

OPENING GAPS IN THE SPECTRUM OF STRICTLY ERGODIC SCHRÖDINGER OPERATORS

ARTUR AVILA^{b,‡}, JAIRO BOCHI*, AND DAVID DAMANIK[‡]

ABSTRACT. We consider Schrödinger operators with dynamically defined potentials arising from continuous sampling along orbits of strictly ergodic transformations. The Gap Labeling Theorem states that the possible gaps in the spectrum can be canonically labelled by an at most countable set defined purely in terms of the dynamics. Which labels actually appear depends on the choice of the sampling function; the missing labels are said to correspond to collapsed gaps. Here we show that for any collapsed gap, the sampling function may be continuously deformed so that the gap immediately opens. As a corollary, we conclude that for generic sampling functions, all gaps are open. The proof is based on the analysis of continuous $SL(2, \mathbb{R})$ cocycles, for which we obtain dynamical results of independent interest.

CONTENTS

1. Introduction	2
2. Statement of Further Results for $SL(2, \mathbb{R})$ -Cocycles	8
3. Outline of the Paper and Discussion of the Methods	10
4. Construction of Almost Invariant Sections	11
5. Construction of Sets With Controlled Return Times	14
6. From Slow Growth to Invariant Sections for the Disk Action	15
7. Solving the Cohomological Equation	17
8. Extending Families of Cocycles: Small Lyapunov Exponents	19
9. Extending Families of Cocycles: Boundedness	24
10. Proof of Theorems 6 and 7	26
Appendix A. Connectedness Considerations	29
Appendix B. Minimal Examples With Many Persistently Closed Gaps	31
Appendix C. Uniform Hyperbolicity and the Rotation Number	32
Appendix D. Proof of the Projection Lemma	35
Appendix E. Mixing Dynamics and Cantor Spectrum	36
References	36

Date: March 12, 2009.

^bCNRS UMR 7599, Laboratoire de Probabilités et Modèles aléatoires. Université Pierre et Marie Curie–Boîte courrier 188. 75252–Paris Cedex 05, France.

[‡]IMPA, Estrada Dona Castorina 110, Jardim Botânico, 22460-320 Rio de Janeiro, Brazil.

*Departamento de Matemática, PUC-Rio, Rua Marquês de São Vicente, 225, 22453-900 Rio de Janeiro, RJ, Brazil.

[‡]Department of Mathematics, Rice University, Houston, TX 77005, USA

We are indebted to Jean Bellissard, Russell Johnson, and Barry Simon for useful comments and hints to the literature. J. B. would like to thank Rice University for its hospitality. D. D. would like to express his gratitude to IMPA and PUC–Rio for their hospitality and financial support through PRONEX–Dynamical Systems and Faperj. This research was partially conducted during the period A. A. served as a Clay Research Fellow. J. B. was partially supported by a grant from CNPq–Brazil. D. D. was supported in part by NSF grants DMS–0653720 and DMS–0800100. We benefited from CNPq and Procad/CAPES support for traveling.

1. INTRODUCTION

1.1. Gap Opening. Let X be a compact metric space and let $f : X \rightarrow X$ be a strictly ergodic homeomorphism (hence all orbits are dense and equidistributed with respect to an invariant probability measure μ). We will also assume that f has a non-periodic finite-dimensional factor.¹

Here we are interested in the one-dimensional discrete Schrödinger operator $H = H_{f,v,x}$ whose potential is obtained by continuous sampling along an orbit of f . Thus H is a bounded self-adjoint operator on $\ell^2(\mathbb{Z})$ given by

$$(1) \quad [H\psi](n) = \psi(n+1) + \psi(n-1) + V(n)\psi(n),$$

where the potential V is given by $V(n) = v(f^n(x))$ for some $x \in X$ and some continuous *sampling function* $v : X \rightarrow \mathbb{R}$.

Since f is minimal, the spectrum of H turns out to depend only on f and v , but not on x , and hence may be denoted by $\Sigma_{f,v}$; it is a perfect compact subset of \mathbb{R} . Strict ergodicity implies that the *integrated density of states* is well defined and only depends on f and v : it is a continuous non-decreasing surjective function $N_{f,v} : \mathbb{R} \rightarrow [0, 1]$. (See §1.5 for details.) The basic relation between $\Sigma_{f,v}$ and $N_{f,v}$ is that $E \notin \Sigma_{f,v}$ if and only if $N_{f,v}$ is constant in a neighborhood of E . Hence, in each bounded connected component of $\mathbb{R} \setminus \Sigma_{f,v}$ (called a *gap* of $\Sigma_{f,v}$), $N_{f,v}$ assumes a constant value $\ell \in (0, 1)$, called the *label* of the gap. A basic question is thus: which values of ℓ appear as labels of gaps?

There is a well known general partial answer to this problem known as the *Gap Labeling Theorem*. Let $G \subset \mathbb{R}$ be the subgroup of all t such that there exist continuous maps $\phi : X \rightarrow \mathbb{R}$ and $\psi : X \rightarrow \mathbb{R}/\mathbb{Z}$ with $t = \int \phi d\mu$ and $\psi(f(x)) - \psi(x) = \phi(x) \bmod 1$. Thus G only depends on f and it is called the *range of the Schwartzman asymptotic cycle*². Often G is a non-discrete group: this happens for instance whenever X has infinitely many connected components.

Gap Labeling Theorem. If ℓ is the label of some gap, then $\ell \in \mathcal{L} = G \cap (0, 1)$.

Of course, the Gap Labeling Theorem does not guarantee that any particular $\ell \in \mathcal{L}$ actually is the label of some gap: it is quite easy to see that the specific choice of the sampling function is also important. A missing label, that is, some $\ell \in \mathcal{L}$ which is not the label of an actual gap, is usually said to correspond to the *collapsed gap* $N_{f,v}^{-1}(\ell)$ (actual gaps are sometimes called *open gaps* for this reason).

Describing which gaps are open in specific situations (and in particular, whether the converse of the Gap Labeling Theorem holds) is usually a very difficult problem. The well known “Dry Ten Martini Problem” (cf. [Si2]) is an instance of this: it asks one to show that all gaps are open for the Almost Mathieu operator ($X = \mathbb{R}/\mathbb{Z}$, $f(x) = x + \alpha$ for some $\alpha \in \mathbb{R} \setminus \mathbb{Q}$ and $v(x) = 2\lambda \cos 2\pi x$ for some $\lambda > 0$). This problem is still open, despite considerable attention (for partial progress, see [CEY], [Pu], [AJ1], and [AJ2]).

Here we are concerned with the more general problem of whether the dynamics, by itself, could *force* the closing of (many) gaps. Questions in this direction arose early in the development of the theory: in a 1982 paper, Belissard asks whether the mixing property would prevent the set of labels of open gaps to be dense in $(0, 1)$ (or equivalently, whether “mixing avoids Cantor spectrum”, [Be1], Open Problem 5).

We will show that the dynamics in fact does not obstruct the opening of gaps, and indeed the converse of the Gap Labeling Theorem holds for generic sampling

¹That is, there is a homeomorphism $g : Y \rightarrow Y$, where Y is a infinite compact subset of some Euclidean space \mathbb{R}^d , and there is a onto continuous map $h : X \rightarrow Y$ such that $h \circ f = g \circ h$.

²See [Sc], [Pa, §5.2], [At]; Theorem 4.2 from [At] justifies this terminology.

functions. This is obtained as a consequence of the following result, which shows that any collapsed gap can be opened by deformation:

Theorem 1 (Gap opening). *For every $v \in C(X, \mathbb{R})$ and any $\ell \in \mathcal{L}$, there exists a continuous path $v_t \in C(X, \mathbb{R})$, $0 \leq t \leq 1$ with $v_0 = v$ such that for every $t > 0$, the spectrum Σ_{f, v_t} has an open gap with label ℓ .*

Since the sampling functions for which the gap with label ℓ is open forms an open subset of $C(X, \mathbb{R})$, we have the following:

Corollary 2 (Generic converse to the Gap Labeling Theorem). *For a generic $v \in C(X, \mathbb{R})$, all gaps are open, that is, for every gap label $\ell \in \mathcal{L}$, the spectrum $\Sigma_{f, v}$ has an open gap with label ℓ .*

In particular, this gives, in our setting, an affirmative answer to another question of Bellissard, who asked in [Be2, Question 4] whether Cantor spectrum is generic whenever it is not precluded by the Gap Labeling Theorem. (A specific mixing example falling in this category is discussed below, answering his earlier question as well.) On the other hand, as we will show in Appendix B, the answer to Bellissard's question becomes negative if one considers slightly more general base dynamics f .

Remark 1.1. We don't know whether for $v \in C(X, \mathbb{R})$ given, it is possible to find a continuous path $v_t \in C(X, \mathbb{R})$, $0 \leq t \leq 1$ with $v_0 = v$ and such that for every $t > 0$ all gaps of Σ_{f, v_t} are open.³

1.2. Examples. One important class of examples (which includes the almost Mathieu operator) is given by almost periodic Schrödinger operators, for which X is an infinite compact Abelian group and f is a minimal translation: in this case \mathcal{L} is a dense subset of $(0, 1)$. Two important subclasses are given by quasi-periodic Schrödinger operators, where $X = \mathbb{R}^d / \mathbb{Z}^d$ for some $d \geq 1$, and by limit-periodic Schrödinger operators, where X is a product of finite cyclic groups.

If \mathcal{L} is not dense in $(0, 1)$, it is impossible for $\Sigma_{f, v}$ to be nowhere dense. On the other hand, when \mathcal{L} is dense in $(0, 1)$, particularly in the almost periodic case, much effort has been dedicated to establishing Cantor spectrum (which means that a dense set of gaps are open). For limit-periodic Schrödinger operators, this issue was studied by various authors in the 1980's (e.g., [M], [AS1], [Ch], [PT]) and, in fact, Corollary 2 is known in this particular case (see [Si1] for an earlier account in the periodic case). However, already for quasi-periodic Schrödinger operators, the problem is much harder.

For the almost Mathieu operator, proving Cantor spectrum corresponds to the Ten Martini Problem, which is now completely solved [AJ1], after several partial results ([BSi], [HS], [CEY], [L], [Pu], [AK]). More generally, the case of non-constant analytic sampling functions along translations of tori has been considered, first by KAM-like methods and more recently by non-perturbative approaches. For Diophantine translations on tori, in any dimension, Eliasson [E] has shown that for sufficiently small generic analytic potentials, all gaps are open. For general non-constant analytic sampling functions v , Goldstein and Schlag [GS] showed that a dense set of gaps can be opened "almost surely" by changing α , in the case where the Lyapunov exponent is positive (which includes the case of large v).

Previous to this work, we had shown in [ABD]⁴ that generically, a dense set of gaps are open, but only for a more restricted set of dynamics ("generalized skew-shifts"): those are the dynamics which fiber over an almost periodic dynamical system. They also have dense \mathcal{L} . An example in this class of particular interest is

³The difficulty is that we deduce Theorem 1 from Theorem 3, and that result does not hold with an extra parameter – see Remark A.4.

⁴Results related to [ABD] can be found in the papers [CF], [FJ], [FJP].

the skew-shift $(x, y) \mapsto (x + \alpha, y + x)$. However, generalized skew-shifts still form a somewhat restricted class; for instance, they are never *weakly mixing*.

In this respect, we would like to mention Fayad's analytic map f on $X = \mathbb{R}^3/\mathbb{Z}^3$ from [Fa] which is strictly ergodic, mixing (hence not a generalized skew-shift), and has dense \mathcal{L} .⁵ Applying our results to this example yields, to the best of our knowledge, the first example of Cantor spectrum for Schrödinger operators arising from mixing dynamics. Of course, an interesting example of a strictly ergodic dynamical system is given by the time one map of the horocycle flow for a compact hyperbolic surface. This is again mixing, but in this case \mathcal{L} is empty!

Other interesting examples of “strictly ergodic dynamics” which are weakly mixing are given by typical interval exchange transformations [AF]. They are not actually strictly ergodic according to the usual definition since they are not homeomorphisms, but this can be bypassed by going to a symbolic setting (blowing up procedure; see [MMY]).

Of course, in the symbolic context, and more generally, among Cantor set maps, there is a wealth of examples, and all have dense \mathcal{L} . See also [BCL] for related examples of homeomorphisms of manifolds.

Remark 1.2. In the analytic setting, it is far from clear whether one should expect that the generic converse of the Gap Labeling Theorem, or even generic Cantor spectrum, remains true for translations on tori in higher (at least 2) dimensions; compare [CS]. We do expect it to be true in one dimension, which is supported by recent advances such as [AJ2] and [GS], but despite this progress, it is not yet known whether the “analytically generic” spectrum is a Cantor set (even restricting to Diophantine translations).

1.3. A Brief History of Gap Labeling. Gap labeling theorems were first established in the early 1980's. In 1982, Johnson and Moser [JM] considered continuum Schrödinger operators, $H = -\frac{d^2}{dx^2} + V$, with almost periodic potential V . For $E \in \mathbb{R}$, if $u = u(x, E)$ is a non-trivial real solution of the associated differential equation $-u'' + Vu = Eu$, then $u'(x, E) + iu(x, E)$ never vanishes; this allowed them to define the rotation number by

$$\rho(E) = \lim_{x \rightarrow \infty} \frac{1}{x} \arg(u'(x, E) + iu(x, E)).$$

They showed that the limit exists, is independent of the solution chosen, and defines a continuous non-decreasing function of E . They also pointed out that due to Sturm oscillation theory, there is a simple relation between the rotation number and the integrated density of states. Since V is almost periodic, $M_V(\lambda) = \lim_{x \rightarrow \infty} \frac{1}{x} \int_0^x V(y) e^{-i\lambda y} dy$ exists for every $\lambda \in \mathbb{R}$ and there are at most countably many λ 's for which $M_V(\lambda) \neq 0$. The set $\mathcal{M}(V)$ of finite integer linear combinations of these λ 's is called the *frequency module* of V . Johnson and Moser then showed that the function ρ is constant precisely on the complement of the spectrum H and, moreover, if $E \in \mathbb{R} \setminus \sigma(H)$, then $2\rho(E) \in \mathcal{M}(V)$. Thus, one can label the gaps of the spectrum by elements of the frequency module of V .

The analogue of the Johnson–Moser results for ergodic discrete Schrödinger operators (and Jacobi matrices) can be found in the 1983 paper [DS] by Delyon and Souillard. In essence, they used Sturm oscillation theory and the easy existence proof for the integrated density of states to define the rotation number in their setting. The relation between the rotation number and the integrated density of states is therefore built into this approach and continues to hold trivially. Moreover, in the case of almost periodic potential, the gaps in the spectrum may again be labeled by elements of the frequency module of the potential.

⁵Details are provided in Appendix E.

Around the same time, Bellissard, Lima, and Testard began to develop an approach to gap labeling based on K -theory of C^* -algebras; see [Be1] for the earliest account and [Be2, Be3, Be4, BLT] for extensive discussions. A major advantage of this approach is that it is not restricted to one space dimension. Let us discuss it in the one-dimensional discrete setting relevant to our present study, but emphasize that the results hold in greater generality. Given a compact space X and a homeomorphism f of X , one can define an associated C^* -algebra $C^*(X, f)$. An f -ergodic probability measure μ induces a trace τ on $C^*(X, f)$. Consider the K_0 group associated with $C^*(X, f)$ and the induced map $\tau_* : K_0(C^*(X, f)) \rightarrow \mathbb{R}$. The abstract gap labeling theorem now reads as follows: Suppose H is a self-adjoint operator affiliated to $C^*(X, f)$ such that its integrated density of states exists with respect to the f -ergodic measure μ . Then the value of the integrated density of states in a gap of the spectrum belongs to $\tau_*(K_0(C^*(X, f))) \cap (0, 1)$. This set of labels is not quite easy to compute, in general. For the case of an irrational rotation of the circle, this was accomplished using the Pimsner–Voiculescu cyclic six-term exact sequence [PV]. This and other computational aspects of the theory are discussed in the references given above. In this context, we would like to also mention the paper [X] by Xia, which shows that the integers in the quasi-periodic gap-labeling can be expressed as Chern characters.

The direct relation between taking the image under the trace homomorphism of the K_0 group of $C^*(X, f)$ and taking the image under Schwartzman’s asymptotic cycle of the first Čech cohomology group $\check{H}^1(Y_f, \mathbb{Z})$, where Y_f is the suspension space of (X, f) , was pointed out by Johnson [J] in 1986; compare also Connes [Co] and the discussions by Bellissard [Be2] and Kellendonk–Zois [KZ]. On the other hand, as pointed out by Bellissard and Scoppola in 1982 [BSc], a naive extension of gap labeling by means of the frequency module to discontinuous sampling function may not work. Indeed, they exhibited an example where the potential arises by coding the rotation of the circle by α irrational with respect to two intervals whose lengths are irrational with respect to α and showed that the integrated density of states takes values not belonging to $\mathbb{Z} + \alpha\mathbb{Z}$.

In recent years, most of the literature on gap labeling has focused on aperiodic tilings and Delone sets in higher dimensions. While this is not immediately related to the subject of this paper, we do want to point out that a recent highlight was accomplished in three simultaneous and independent works by Bellissard–Benedetti–Gambaudo [BBG], Benameur–Oyono–Oyono [BO], and Kaminker–Putnam [KP]. They proved a version of the gap labeling theorem for aperiodic, repetitive tilings or Delone sets of finite local complexity.

1.4. Cocycle Dynamics. It is well established that spectral properties of one-dimensional Schrödinger operators are intimately connected to a study of the solutions to the time-independent Schrödinger equation. This equation in turn can be reformulated in terms of transfer matrices. Since the potential is generated by sampling along the orbits of f and the dependence of the transfer matrices on the potential is local, it turns out that these matrices have a so-called cocycle structure. In particular, the dynamics of these $\text{SL}(2, \mathbb{R})$ -valued cocycles describe the behavior of solutions and hence are closely related to spectral information. We will make this connection more explicit in the next subsection. Here, we discuss basic notions related to $\text{SL}(2, \mathbb{R})$ -valued cocycles and present a relative of Theorem 1 in this context.

Given a continuous map $A : X \rightarrow \text{SL}(2, \mathbb{R})$, we consider the skew-product $X \times \text{SL}(2, \mathbb{R}) \rightarrow X \times \text{SL}(2, \mathbb{R})$ given by $(x, g) \mapsto (f(x), A(x) \cdot g)$. This map is called the

cocycle (f, A) .⁶ For $n \in \mathbb{Z}$, A^n is defined by $(f, A)^n = (f^n, A^n)$. The *Lyapunov exponent* of the cocycle is the non-negative number:

$$L(A) = \lim_{n \rightarrow \infty} \frac{1}{n} \log \|A^n(x)\|.$$

(The limit exists for μ -almost every x .)

We say that a cocycle (f, A) is *uniformly hyperbolic* if there exist constants $c > 0$, $\lambda > 1$ such that $\|A^n(x)\| > c\lambda^n$ for every $x \in X$ and $n > 0$.⁷ This is equivalent to the usual hyperbolic splitting condition: see [Y]. Thus uniform hyperbolicity is an open condition in $C(X, \text{SL}(2, \mathbb{R}))$. Denote

$$\mathcal{UH} = \{A \in C(X, \text{SL}(2, \mathbb{R})); (f, A) \text{ is uniformly hyperbolic}\}.$$

We say that two cocycles (f, A) and (f, \tilde{A}) are *conjugate* if there exists a *conjugacy* $B \in C(X, \text{SL}(2, \mathbb{R}))$ such that $\tilde{A}(x) = B(f(x))A(x)B(x)^{-1}$. Uniform hyperbolicity is invariant under conjugation.

Next suppose $A \in C(X, \text{SL}(2, \mathbb{R}))$ is homotopic to a constant. Then one can define a homeomorphism $F : X \times \mathbb{R} \rightarrow X \times \mathbb{R}$ so that $F(x, t) = (f(x), F_2(x, t))$ and $e^{2\pi i F_2(x, t)}$ is a positive real multiple of $A(x) \cdot e^{2\pi i t}$, with the usual identification of \mathbb{C} with \mathbb{R}^2 . We call $F_2(x, t)$ a lift of A . The lift is not uniquely defined, since for any continuous function $\phi : X \rightarrow \mathbb{Z}$, $(x, t) \mapsto F_2(x, t) + \phi(x)$ is also a lift, but this is the only possible ambiguity.

Write $F_2^n(x, t)$ for the second component of $F^n(x, t)$. Then the limit

$$\rho_{f, A} = \lim_{n \rightarrow \infty} \frac{F_2^n(x, t) - t}{n}$$

exists uniformly and is independent of (x, t) , it is called the *fibered rotation number* of the cocycle (f, A) ; compare [H]. Due to the ambiguity in the choice of F_2 , ρ is only defined up to addition of an element of the group $G' = \{\int_X \phi \, d\mu : \phi \in C(X, \mathbb{Z})\}$. Clearly, $\mathbb{Z} \subset G' \subset G$. Moreover, if X is connected, then $G' = \mathbb{Z}$, while if it is totally disconnected, then $G' = G$.

We will be interested in the variation of the fibered rotation number. For any map $A_0 : X \rightarrow \text{SL}(2, \mathbb{R})$ that is homotopic to a constant, there is a neighborhood \mathcal{V} of it in $C(X, \text{SL}(2, \mathbb{R}))$ such that for each $A \in \mathcal{V}$, we can continuously select a lift $F_A : X \times \mathbb{R} \rightarrow X \times \mathbb{R}$. Using those lifts, we define a map $A \in \mathcal{V} \mapsto \rho_{f, A} \in \mathbb{R}$, which is in fact continuous (see [H], §5.7), and is called a *continuous determination of the fibered rotation number*.

It turns out that (see [J])

$$(2) \quad A \in \mathcal{UH} \Rightarrow 2\rho_{f, A} \in G.$$

We will establish the following variant of Theorem 1 for $\text{SL}(2, \mathbb{R})$ -cocycles:

Theorem 3 (Accessibility by uniformly hyperbolic cocycles). *Suppose $A \in C(X, \text{SL}(2, \mathbb{R}))$ is homotopic to a constant and obeys $2\rho_{f, A} \in G$. Then there exists a continuous path $A_t \in C(X, \text{SL}(2, \mathbb{R}))$, $0 \leq t \leq 1$ with $A_0 = A$ and such that $A_t \in \mathcal{UH}$ for every $t > 0$. (In particular, ρ_{f, A_t} is independent of t .)*

In fact, as we will see in the next subsection, Theorem 1 may be derived from Theorem 3.

⁶Since f is fixed, we sometimes call A a cocycle.

⁷If A is a real 2×2 matrix then $\|A\| = \sup_{\|v\| \neq 0} \|A(v)\|/\|v\|$, where $\|v\|$ is the Euclidean norm of $v \in \mathbb{R}^2$.

1.5. Schrödinger Operators and Relation to Dynamics. Let us now review in more detail the basic objects introduced in the beginning and their relation to the dynamics of cocycles.

Consider the operator $H = H_{f,v,x}$ given by (1). Since f is minimal and v is continuous, the x -independence of the spectrum $\Sigma_{f,v}$ of H can be obtained from strong operator convergence modulo conjugation with shifts. This common spectrum has a convenient description in terms of the integrated density of states, which is defined as follows. There is a probability measure $dN_{f,v}$ that obeys

$$\int \langle \delta_0, g(H) \delta_0 \rangle d\mu(x) = \int g(E) dN_{f,v}(E)$$

for every bounded, measurable function g . Indeed, by the spectral theorem, $dN_{f,v}$ is given by the μ -average of the x -dependent spectral measures associated with H and δ_0 . The distribution function of $dN_{f,v}$ is called the integrated density of states and denoted by $N_{f,v}$. By construction, it is a non-decreasing function from \mathbb{R} onto $[0, 1]$, which is constant in every connected component of $\mathbb{R} \setminus \Sigma_{f,v}$. Moreover, it is not hard to show that it is continuous and strictly increasing on $\Sigma_{f,v}$. That is, the measure $dN_{f,v}$ is non-atomic and its topological support is equal to $\Sigma_{f,v}$. There is an alternate way to define $N_{f,v}(E)$ by restricting H to finite intervals, counting the eigenvalues below E , and taking a thermodynamic limit; compare [AS2]. A third way is explained below.

Spectral properties of general Schrödinger operators of the form (1) are most conveniently studied in terms of the solutions to the one-parameter family of difference equations

$$(3) \quad u(n+1) + u(n-1) + V(n)u(n) = Eu(n),$$

where the energy E belongs to \mathbb{R} . Equivalently,

$$\begin{pmatrix} u(n+1) \\ u(n) \end{pmatrix} = \begin{pmatrix} E - V(n) & -1 \\ 1 & 0 \end{pmatrix} \begin{pmatrix} u(n) \\ u(n-1) \end{pmatrix}.$$

Since in our situation, we have $V(n) = v(f^n x)$, we may consider the cocycle $(f, A_{E,v})$, where

$$A_{E,v}(x) = \begin{pmatrix} E - v(x) & -1 \\ 1 & 0 \end{pmatrix},$$

and observe that u solves (3) if and only if

$$\begin{pmatrix} u(n) \\ u(n-1) \end{pmatrix} = A_{E,v}^n \begin{pmatrix} u(0) \\ u(-1) \end{pmatrix}.$$

Consequently, (f, A) is called a *Schrödinger cocycle* if A takes its values in the set

$$(4) \quad \mathcal{S} = \left\{ \begin{pmatrix} t & -1 \\ 1 & 0 \end{pmatrix}; t \in \mathbb{R} \right\} \subset \text{SL}(2, \mathbb{R}).$$

One of the fundamental results linking the spectral theory of H and the dynamics of the family of Schrödinger cocycles $(f, A_{E,v})$ is that $\Sigma_{f,v}$ consists of those energies E for which $(f, A_{E,v})$ is not uniformly hyperbolic:

$$(5) \quad \mathbb{R} \setminus \Sigma_{f,v} = \{E \in \mathbb{R} : (f, A_{E,v}) \text{ is uniformly hyperbolic}\};$$

see Johnson [J].

The maps $A_{E,v}$ are homotopic to a constant and hence the cocycle $(f, A_{E,v})$ has an associated rotation number which we denote by $\rho_{f,v}(E)$. For Schrödinger cocycles, there is a canonical way to remove the ambiguity in the definition of ρ ,⁸ so that it can be interpreted as an element of $[0, 1/2]$ (and not of \mathbb{R}/G' as in the

⁸Namely, one can uniquely choose the lift F_2 so that $F_2(x, 1/4) = 1/2$ for every $x \in X$.

general case). There is a simple relation between $N_{f,v}$ and $\rho_{f,v}$ which follows from Sturm oscillation theory:

$$(6) \quad N_{f,v}(E) = 1 - 2\rho_{f,v}(E).$$

It follows from (2), (5), and (6) that the gaps of $\Sigma_{f,v}$ (i.e., the bounded connected components of $\mathbb{R} \setminus \Sigma_{f,v}$) may be labeled by elements of $\mathcal{L} = G \cap (0, 1)$ so that the label of a gap corresponds to the constant value $N_{f,v}$ takes on it.

To deduce Theorem 1 from Theorem 3, we use the following result, which is an improved version of [ABD, Lemma 9] and is proved in Appendix D. (Recall that $\mathcal{S} \subset \text{SL}(2, \mathbb{R})$ indicates the set of Schrödinger matrices, as in (4).)

Lemma 1.3 (Projection Lemma). *Let $f : X \rightarrow X$ be a minimal homeomorphism of a compact metric space with at least 4 points, and let $A \in C^0(X, \mathcal{S})$. Then there exist a neighborhood $\mathcal{W} \subset C^0(X, \text{SL}(2, \mathbb{R}))$ of A and continuous maps*

$$\Phi = \Phi_A : \mathcal{W} \rightarrow C^0(X, \mathcal{S}) \quad \text{and} \quad \Psi = \Psi_A : \mathcal{W} \rightarrow C^0(X, \text{SL}(2, \mathbb{R}))$$

satisfying:

$$\begin{aligned} \Psi(B)(f(x)) \cdot B(x) \cdot [\Psi(B)(x)]^{-1} &= \Phi(B)(x), \\ \Phi(A) &= A \text{ and } \Psi(A) = \text{id}. \end{aligned}$$

In particular, an $\text{SL}(2, \mathbb{R})$ -valued perturbation of a Schrödinger cocycle is conjugate to an \mathcal{S} -valued perturbation.

Proof of Theorem 1. Take a sampling function $v \in C(X, \mathbb{R})$ and a label $\ell \in \mathcal{L} = G \cap (0, 1)$. Assume the gap of $\Sigma_{f,v}$ of label ℓ is collapsed, that is, there is a unique $E_0 \in \mathbb{R}$ such that $N_{f,v}(E_0) = \ell$. Let

$$A(x) = \begin{pmatrix} E_0 - v(x) & -1 \\ 1 & 0 \end{pmatrix}.$$

By (6), $\rho(A) = (1 - \ell)/2$. So we can apply Theorem 3 and find a continuous family of cocycles A_t , $t \in [0, 1]$ with $A_0 = A$ and $A_t \in \mathcal{UH}$ for $t > 0$.

Apply Lemma 1.3 and define $\tilde{A}_t = \Phi_A(A_t)$ for small t . Since $\tilde{A}_t \in C(X, \mathcal{S})$, there is a continuous path $v_t \in C(X, \mathbb{R})$ so that

$$\tilde{A}_t(x) = \begin{pmatrix} E_0 - v_t(x) & -1 \\ 1 & 0 \end{pmatrix}.$$

For $t > 0$, we have $\tilde{A}_t \in \mathcal{UH}$ and so, by (5), $E_0 \notin \Sigma_{f,v_t}$. Thus $N_{f,v_t}(E_0)$, which equals $1 - 2\rho(\tilde{A}_t) = \ell$, is a label of an open gap. \square

2. STATEMENT OF FURTHER RESULTS FOR $\text{SL}(2, \mathbb{R})$ -COCYCLES

This paper actually contains more refined results on cocycles, and Theorem 3 is obtained as a corollary of one of them, as we now explain.

A cocycle (f, A) is called *bounded* if there exists $C > 1$ such that $\|A^n(x)\| \leq C$ for every $x \in X$ and $n \in \mathbb{Z}$. It is known that a cocycle is bounded if and only if it is conjugate to a cocycle of rotations, that is, to (f, \tilde{A}) with $\tilde{A} : X \rightarrow \text{SO}(2, \mathbb{R})$; see [Ca], [EJ], [Y].

Our next main result is:

Theorem 4 (Accessibility by bounded cocycles). *If $A \in C(X, \text{SL}(2, \mathbb{R})) \setminus \mathcal{UH}$, then there exists a continuous path $A_t \in C(X, \text{SL}(2, \mathbb{R}))$, $t \in [0, 1]$, such that $A_0 = A$ and (f, A_t) is conjugate to an $\text{SO}(2, \mathbb{R})$ -valued cocycle for every $t \in (0, 1]$. Moreover, the conjugacy can be chosen to depend continuously on $t \in (0, 1]$.*

In particular, bounded cocycles are dense in the complement of \mathcal{UH} . In fact, we are able to prove this fact directly and with weaker hypotheses on the dynamics:

Theorem 5 (Denseness of bounded cocycles). *Assume that f is a minimal homeomorphism of a compact metric space X such that at least one of the following holds:*

- a) X is finite-dimensional;
- b) f has at most countably many distinct ergodic invariant measures, and f has a non-periodic finite-dimensional factor.

Then, the bounded cocycles form a dense subset of $C(X, \mathrm{SL}(2, \mathbb{R})) \setminus \mathcal{UH}$.

We now focus on the case of cocycles that are homotopic to a constant, which form a subset of $C(X, \mathrm{SL}(2, \mathbb{R}))$ that we indicate by $C_0(X, \mathrm{SL}(2, \mathbb{R}))$.

Given $A \in C_0(X, \mathrm{SL}(2, \mathbb{R}))$, fix a continuous determination ρ of the fibered rotation number on a neighborhood of it. We say that the cocycle A is *locked* if ρ is constant on a neighborhood of A . We say that A is *semi-locked* if it is not locked and it is a point of local extremum for the function ρ . Otherwise we say that A is *unlocked*. This classification is obviously independent of the chosen determination of ρ .

In fact, a cocycle is locked if and only if it is uniformly hyperbolic, and any semi-locked cocycle A can be accessed by uniformly hyperbolic cocycles. The proof of these statements and some related facts may be found in Appendix C. (Similar results for “non-linear cocycles” over rotations can be found in [BJ].)

For unlocked cocycles, we can strengthen the conclusion of Theorem 4 considerably:

Theorem 6 (Accessibility by cocycles conjugate to a rotation). *Let $A \in C_0(X, \mathrm{SL}(2, \mathbb{R}))$ be an unlocked cocycle. Let $\alpha \in \mathbb{R}$ be such that $\alpha = \rho_{f,A} \bmod G$. Then there exist continuous paths $A_t \in C(X, \mathrm{SL}(2, \mathbb{R}))$, $t \in [0, 1]$, and $B_t \in C(X, \mathrm{SL}(2, \mathbb{R}))$, $t \in (0, 1]$, such that $A_0 = A$ and $B_t(f(x))^{-1}A_t(x)B_t(x) = R_{2\pi\alpha}$ for $t \in (0, 1]$.*

The assumption that the cocycle be unlocked is necessary: see Remark 10.6.

It is useful to consider also a weaker notion of conjugacy: We say that two cocycles (f, A) and (f, \tilde{A}) are *projectively conjugate* if there exists a *projective conjugacy* $B \in C(X, \mathrm{PSL}(2, \mathbb{R}))$ such that $\tilde{A}(x) = B(f(x))A(x)B(x)^{-1}$ modulo sign. The corresponding version of Theorem 6 is:

Theorem 7 (Accessibility by cocycles projectively conjugate to a rotation). *Let $A \in C_0(X, \mathrm{SL}(2, \mathbb{R}))$ be an unlocked cocycle. Let $\alpha \in \mathbb{R}$ be such that $2\alpha = 2\rho_{f,A} \bmod G$. Then there exist continuous paths $A_t \in C(X, \mathrm{SL}(2, \mathbb{R}))$, $t \in [0, 1]$, and $B_t \in C(X, \mathrm{PSL}(2, \mathbb{R}))$, $t \in (0, 1]$, such that $A_0 = A$ and $B_t(f(x))^{-1}A_t(x)B_t(x) = R_{2\pi\alpha}$ (modulo sign) for $t \in (0, 1]$.*

This result easily yields Theorem 3:

Proof of Theorem 3. Suppose $A \in C(X, \mathrm{SL}(2, \mathbb{R}))$ is homotopic to a constant with $2\rho_{f,A} \in G$; we want to access A by a path A_t in \mathcal{UH} . As mentioned above, if A is locked, then A is in \mathcal{UH} and thus we can take a constant path. If A is semi-locked, then a path of the form $A_t = R_{\varepsilon t}A$ (for some small $\varepsilon \neq 0$) works; see Proposition C.1. Thus assume that A is unlocked.

Apply Theorem 7 with $\alpha = 0$ and find continuous paths $A_t \in C(X, \mathrm{SL}(2, \mathbb{R}))$, $t \in [0, 1]$, and $B_t \in C(X, \mathrm{PSL}(2, \mathbb{R}))$, $t \in (0, 1]$, such that $A_0 = A$ and $A_t(x) = B_t(f(x))B_t(x)^{-1}$ (modulo sign) for $t \in (0, 1]$. Let $\tau : (0, 1] \rightarrow (1, \infty)$ be a continuous

function. Take $\tilde{A}_t \in C(X, \mathrm{SL}(2, \mathbb{R}))$ such that

$$\tilde{A}_t(x) = B_t(f(x)) \begin{pmatrix} \tau(t) & 0 \\ 0 & 1/\tau(t) \end{pmatrix} B_t(x)^{-1} \pmod{\text{sign}}.$$

Choosing a function $\tau(t)$ that goes sufficiently fast to 1 as $t \rightarrow 0$, we have $\lim_{t \rightarrow 0} \tilde{A}_t(x) = A(x)$. This path has the desired properties. \square

As for uniformly hyperbolic cocycles, they cannot, in general, be approximated by cocycles that are conjugate to a constant (see Remark 3 in [ABD]). But for projective conjugacy there is no obstruction; indeed we show:

Theorem 8 (Accessibility by cocycles projectively conjugate to a hyperbolic matrix). *Let $A \in \mathcal{UH}$, and let $D \in \mathrm{SL}(2, \mathbb{R})$ be such that $L(A) = L(D)$. Then there exist continuous paths $A_t \in C(X, \mathrm{SL}(2, \mathbb{R}))$, $t \in [0, 1]$, and $B_t \in C(X, \mathrm{PSL}(2, \mathbb{R}))$, $t \in (0, 1]$, such that $A_0 = A$ and $B_t(f(x))A_t(x)B_t(x)^{-1} = D$ (in $\mathrm{PSL}(2, \mathbb{R})$) for $t \in (0, 1]$.*

Remark 2.1. Cocycles that are projectively conjugate to a constant are called *reducible*. Theorems 7 and 8 (together with the fact that the set of semi-locked cocycles has empty interior) imply that reducibility is dense in $C_0(X, \mathrm{SL}(2, \mathbb{R}))$. As for uniform hyperbolicity, it is dense in $C_0(X, \mathrm{SL}(2, \mathbb{R}))$ if and only if G is a non-discrete subgroup of \mathbb{R} , by Theorem 3 (and the basic fact (2)). This statement cannot be significantly improved; compare Theorem 9 in Appendix B.

More generally, as in [ABD], we can let *Ruth* be the set of all cocycles that are reducible up to homotopy (or, equivalently, are projectively conjugate to a cocycle $C_0(X, \mathrm{SL}(2, \mathbb{R}))$; see Lemma 8 in [ABD]). Any uniformly hyperbolic cocycle is reducible up to homotopy, and in [ABD] it is shown that reducible uniformly hyperbolic cocycles are dense in *Ruth* in the particular case where f is a generalized skew-product. Our results imply that in our more general context, reducibility is still dense in *Ruth*, and uniformly hyperbolic cocycles are dense in *Ruth* if and only if G is non-discrete.

Remark 2.2. We also notice the following consequence of our methods: any continuous cocycle A that is not uniformly hyperbolic admits arbitrarily small neighborhoods in $C(X, \mathrm{SL}(2, \mathbb{R})) \setminus \mathcal{UH}$ that are path connected (and even have trivial homotopy groups in all dimensions). See the comment after Proposition 9.2.

3. OUTLINE OF THE PAPER AND DISCUSSION OF THE METHODS

Most of the paper is dedicated to building up the tools that will be involved in the proof of cousin Theorems 6 and 7, with the other results being obtained on the side. These tools can be split into two classes of results which produce families of cocycles with certain features. The input in either case is always some initial family of cocycles, but while the first kind of result is only concerned with outputting a suitable *perturbation* of the family, the second is primarily concerned with the *extension* of the parameter space.

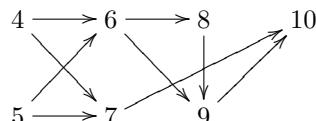
The perturbation arguments are developed in Sections 4 to 7: they refine the ideas introduced in [ABD] to deal with individual cocycles over generalized skew-shifts, both to address a more general class of dynamical systems and to allow us to work in the context of parametrized families. Since many problems about cocycles can be rephrased in terms of the existence of invariant sections for an associated skew-product (for instance, conjugacy to rotations is equivalent to the existence of an invariant section for the associated skew-product arising from the disk action of $\mathrm{SL}(2, \mathbb{R})$), we first (Section 4) discuss in a more abstract setting the problem of existence of *almost invariant sections*, satisfying the invariance condition except in some specific set. The actual estimates for the section depend on the recurrence

properties of the set and the existence of sets with good recurrence properties is the topic of Section 5. Those results are used in Section 6 to show that a family of $\text{SL}(2, \mathbb{R})$ -cocycles with small Lyapunov exponents can be perturbed to become continuously conjugate to a family of cocycles of rotations (Proposition 6.3); hence Theorem 5 follows as a particular case. A similar result (Lemma 7.1), regarding the cohomological equation (which plays a role when conjugating a cocycle of rotations to a constant), is obtained in Section 7, and permits us to establish Theorem 8.

The extension arguments are developed in Sections 8 and 9: the questions considered are rather novel, and their proofs need the introduction of several ideas. In Section 8 we show that a family of cocycles with zero Lyapunov exponent, defined over the boundary of a cube, can be extended to a family of cocycles with small Lyapunov exponents, defined over the whole cube, provided there is no topological obstruction (Lemma 8.1). The proof is technically quite involved, jumping from the continuous to the measurable category and back, and is mostly independent from the rest of the paper (which never departs from the continuous category): it is also the only place where ergodic theory plays a significant role. In Section 9, Lemma 8.1 is used to bootstrap a considerably more refined result (Proposition 9.2): under the weaker condition that the boundary values are not uniformly hyperbolic, one gets the stronger conclusion that the extended family can be continuously conjugated to cocycles of rotations in the interior of the cube. In particular, we obtain Theorem 4.

In Section 10, everything is put together in the proof of Theorems 6 and 7. We first show (Lemma 10.3) that given A unlocked, it is possible to construct a two-parameter family $A_{s,t}$, $(s,t) \in [0,1]^2$ with $A_{0,0} = A$, such that $\rho(A_{0,s}) < \rho(A) < \rho(A_{s,0})$, and which is continuously conjugate to rotations (except possibly at $(0,0)$). This almost implies that A is accessible from the interior of the square $[0,1]^2$ through a path with constant rotation number. The actual accessibility of A by cocycles conjugate to a fixed rotation is obtained with a little wiggling involving in particular the solution of the cohomological equation.

The following diagram gives the logical dependence between the sections:



The paper also contains five appendices. The first contains complementary results and examples concerning connectedness of some sets of cocycles. The second shows that complete gap labeling, and even just the Cantor structure of the spectrum, ceases to be generic when the assumption of unique ergodicity is dropped. The third exposes some facts on locked / semi-locked / unlocked cocycles. The fourth proves the projection lemma that was used to realize $\text{SL}(2, \mathbb{R})$ perturbations by Schrödinger matrices. Finally, in the last appendix we show how Fayad's results can be applied to obtain a mixing strictly ergodic example with a dense set of labels.

4. CONSTRUCTION OF ALMOST INVARIANT SECTIONS

In this section and in the next section, $f : X \rightarrow X$ is an arbitrary homeomorphism of a compact metric space. In fact, even if we are ultimately interested in minimal dynamical systems, in order to obtain “parameterized” results in the later sections, we will apply the results obtained here to maps of the form $\hat{f}(x, t) = (f(x), t)$ (the second coordinate representing the parameter).

Let Y be a topological space and let $F : X \times Y \rightarrow X \times Y$ be an invertible skew-product over f , that is, a homeomorphism of the form $(x, y) \mapsto (f(x), F_x(y))$.

An *invariant section* for F is a continuous map $y : X \rightarrow Y$ such that $F_x(y(x)) = y(f(x))$. Throughout the paper, a common theme will be the construction of *almost invariant sections*, which satisfy $F_x(y(x)) = y(f(x))$ in much of X . We define the *support* of an almost invariant section as the set of all x such that $F_x(y(x)) \neq y(f(x))$. We will also be interested in estimating “how large” such almost invariant sections are. In the abstract setting considered here, this largeness will be measured by a continuous function $M : Y \rightarrow [0, \infty)$.

We say that the pair (Y, M) has the *Tietze property* if for any metric space L , every continuous function $y : K \rightarrow Y$ defined on a compact subset $K \subset L$ can be extended to a continuous map $y : L \rightarrow Y$ such that

$$\sup_{x \in L} M(y(x)) = \sup_{x \in K} M(y(x)).$$

The following are the examples of pairs (Y, M) that we will use:

- $Y = \mathbb{R}$ and $M(y) = |y|$;
- Y is the unit disk \mathbb{D} and $M(y) = d(y, 0)$, where d is the hyperbolic distance (see §6.1).

Using Tietze’s Extension Theorem, we see that the Tietze property holds in these two examples.

We denote

$$\|F\|_M = \sup_{(x,y) \in X \times Y} |M(F_x(y)) - M(y)| \in [0, \infty].$$

Notice that $\|F^{-1}\|_M = \|F\|_M$.

To construct and estimate almost invariant sections whose support is contained in the interior of a compact set K , we will need only a few aspects of the recurrence of the dynamics to K , which we encode in the following notions. We say that $K \subset X$ is *n-good* if its first n iterates are disjoint, that is, $f^i(K) \cap K = \emptyset$ for $1 \leq i \leq n-1$. Let us say that $K \subset X$ is *N-spanning* if the first N iterates cover X , that is, $\bigcup_{i=0}^{N-1} f^i(K) = X$. We say that $K \subset X$ is *d-mild* if the orbit $\{f^n x : n \in \mathbb{Z}\}$ of each point in X enters ∂K at most d times. The construction of sets with appropriate recurrence properties, under suitable dynamical assumptions, will be carried out in the next section.

The following abstract lemma is the main result of this section. (A simple particular case of it, where the dynamics fibers over an irrational rotation, corresponds to Lemma 2 from [ABD].)

Lemma 4.1. *Let $f : X \rightarrow X$ be a homeomorphism of a compact metric space X , and let $K \subset X$ be a d -mild, n -good, N -spanning set. Let $\Lambda \subset X$ be a compact (possibly empty) f -invariant set. Let (Y, M) be a pair with Tietze property, and let $F : X \times Y \rightarrow X \times Y$ be an invertible skew-product over f . Let $y_0 : X \rightarrow Y$ be an almost invariant section supported outside Λ . Then there exists an almost invariant section $y : X \rightarrow Y$ supported in $\text{int } K$ with the following properties:*

- a) y equals y_0 on Λ ;
- b) $\sup_{x \in X} M(y(x)) \leq \sup_{x \in X} M(y_0(x)) + d \max_{j \in [n, N]} \|F^j\|_M$.

Proof. We begin by using the set K and the dynamics f to decompose the space X .

For each $x \in X$, let

$$\begin{aligned} \ell^+(x) &= \min\{j \geq 0 : f^j(x) \in \text{int } K\}, \quad \text{and} \quad \ell^-(x) = \min\{j > 0 : f^{-j}(x) \in \text{int } K\}; \\ T(x) &= \{j \in \mathbb{Z} : -\ell^-(x) < j < \ell^+(x)\} \quad \text{and} \quad T_B(x) = \{j \in T(x) : f^j(x) \in \partial K\}. \end{aligned}$$

Notice T and T_B are upper-semicontinuous. Define $N(x) = \#T_B(x)$ and $X^i = \{x \in X : N(x) \geq d - i\}$. The sets X^i are closed and

$$X = X^d \supset X^{d-1} \supset \cdots \supset X^0 \supset X^{-1} = \emptyset.$$

Let also $Z^i = X^i \setminus X^{i-1} = \{x : N(x) = d - i\}$.

Claim. ℓ^+ is locally constant on the set Z^i .

Proof. We will show that T_B and T are locally constant on Z^i . Fix $x \in Z^i$ and let $y \in Z^i$ be very close to x . We have $T_B(y) \subset T_B(x)$. Since y is also in Z^i , we must have equality. Now, if $j \in T(x)$, then either $f^j(x) \in \partial K$ or $f^j(x) \in X \setminus K$. In the former case, $f^j(y) \in \partial K$ (because $T_B(y) = T_B(x)$), and in the latter, $f^j(y)$ is also in $X \setminus K$ (because y is close to x). Since $f^{\ell^+(x)}(y), f^{-\ell^-(x)}(y) \notin \text{int } K$, we conclude that $T(y) = T(x)$. \square

Let $Z_\ell^i = \{x \in Z^i : \ell^+(x) = \ell\}$; by the claim this is a (relatively) open subset of X^i . Also notice the following facts:

$$\begin{aligned} X^i &= X^{i-1} \sqcup Z_0^i \sqcup Z_1^i \sqcup \cdots \sqcup Z_N^i, & \overline{Z_\ell^i} \setminus Z_\ell^i &\subset X^{i-1}, \\ Z_0^i &= Z^i \cap \text{int } K, & f(Z_\ell^i) &= Z_{\ell-1}^i \text{ if } \ell > 0. \end{aligned}$$

Now let (Y, M) , F , Λ , y_0 be as in the assumptions of the lemma. We will describe a procedure to successively define y on the set $\Lambda \cup (X^0 \cap \text{int } K)$, then on $\Lambda \cup X^0$, then on $\Lambda \cup X^0 \cup (X^1 \cap \text{int } K)$, then on $\Lambda \cup X^1$, then on $\Lambda \cup X^1 \cup (X^2 \cap \text{int } K)$, \dots , then on $\Lambda \cup X^{d-1} \cup (X^d \cap \text{int } K)$, and finally on $\Lambda \cup X^d = X$.

We begin by defining y on $\Lambda \cup (X^0 \cap \text{int } K)$ as equal to y_0 .

Let $i \in \{0, \dots, d\}$. Assume, by induction, that the map y has already been continuously defined on the set $\Lambda \cup X^{i-1} \cup (X^i \cap \text{int } K)$, and has the two properties below:

(7) $F(x, y(x)) = (f(x), y(f(x)))$ for every x in the domain and not in $\text{int } K$;

$$(8) \quad \sup_{x \in \Lambda \cup X^{i-1} \cup (X^i \cap \text{int } K)} M(y(x)) \leq \sup_{x \in X} M(y_0(x)) + iA,$$

where $A = \max_{j \in [n, N]} \|F^j\|_M$.

Now, for each $\ell = 1, \dots, N$, the map y is already defined on $f^\ell(Z_\ell^i) \subset X^i \cap \text{int } K$. Naturally, we define y on Z_ℓ^i by

$$y(x) = (F_x^\ell)^{-1} (y(f^\ell(x))).$$

This defines y on the set $\Lambda \cup X^i$ in such a way that (7) holds, and moreover

$$\sup_{x \in \Lambda \cup X^i} M(y(x)) \leq \sup_{x \in \Lambda \cup X^{i-1} \cup (X^i \cap \text{int } K)} M(y(x)) + A.$$

Let us check that the map y so defined on $\Lambda \cup X^i$ is continuous. Since Λ and X^i are closed sets and $y|_\Lambda = y_0|_\Lambda$ is continuous, we only need to check that $y|_{X^i}$ is continuous. Take a sequence x_j in X^i converging to some x and let us show that $y(x_j) \rightarrow y(x)$. If $x \in Z^i$, convergence follows from the claim above and the continuity of $y|_{X^i \cap \text{int } K}$. Next assume $x \in X^{i-1}$. Since $X^i = X^{i-1} \sqcup Z_0^i \sqcup Z_1^i \sqcup \cdots \sqcup Z_N^i$, and $y|_{X^{i-1}}$ is continuous, it suffices to consider the case where the sequence x_j is contained in some Z_ℓ^i . Then $f^\ell(x_j)$ and $f^\ell(x)$ all belong to $(X^i \cap \text{int } K) \cup X^{i-1}$, where y is continuous. So, applying $F^{-\ell}$, we obtain $y(x) = \lim y(x_j)$, as desired.

At this point, y is continuously defined on $\Lambda \cup X^i$. If $i = d$, we are done. Otherwise, we use the Tietze property to extend continuously y to $\Lambda \cup X^i \cup (X^{i+1} \cap \text{int } K)$ in a way such that

$$\sup_{x \in \Lambda \cup X^i \cup (X^{i+1} \cap \text{int } K)} M(y(x)) = \sup_{x \in \Lambda \cup X^i} M(y(x)).$$

Then the hypotheses (7) and (8) hold with $i+1$ in place of i ; thus we can increment i and continue the procedure. \square

5. CONSTRUCTION OF SETS WITH CONTROLLED RETURN TIMES

In the applications of Lemma 4.1, it would be bad if N was much bigger than n (because then we wouldn't get a useful estimate for the M -size of the almost invariant section). This issue is settled by means of the following result:

Proposition 5.1. *Let $f : X \rightarrow X$ be a homeomorphism of a compact metric space that admits a non-periodic minimal finite-dimensional factor. Then there is $d \in \mathbb{N}$ such that for every $n \in \mathbb{N}$, there is a compact set $K \subset X$ that is n -good, $(d+2)n-1$ -spanning, and d -mild.*

We mention that in the case where f fibers over an irrational rotation, the proposition can be proven (with $d=1$) by a simple explicit construction based on continued fractions – see Figure 1 from [ABD].

If $\tilde{f} : \tilde{X} \rightarrow \tilde{X}$ is a factor of $f : X \rightarrow X$ and $\tilde{K} \subset \tilde{X}$ is d -mild, n -good and N -spanning for \tilde{f} , then $\pi^{-1}(\tilde{K})$ is d -mild, n -good and N -spanning for f , where $\pi : X \rightarrow \tilde{X}$ is a continuous surjective map such that $\pi \circ f = \tilde{f} \circ \pi$. Thus it suffices to prove this statement for the minimal non-periodic finite-dimensional factor of f . For the remainder of this section, it will therefore be assumed that X is finite-dimensional and infinite and that f is minimal.

Lemma 5.2. *There exists $d \in \mathbb{N}$ such that for every compact set $K \subset X$ and every neighborhood U of K , there exists a compact neighborhood $K' \subset U$ of K that is d -mild.*

Proof. This follows easily from Lemmas 4 and 5 in [AB]. \square

Let now d given by Lemma 5.2 be fixed.

Lemma 5.3. *Let $K \subset X$ be an n -good, N -spanning compact set. Then there exists a compact neighborhood K' of K that is n -good, N -spanning and d -mild.*

Proof. If K' is a compact neighborhood of K sufficiently close to K , then it is still n -good and N -spanning. So the result follows from Lemma 5.2. \square

Lemma 5.4. *For every $n \in \mathbb{N}$, there exists an n -good, $(d+2)n-1$ -spanning compact set $K \subset X$.*

Proof. Let n be fixed. Let us start with a compact set K with non-empty interior whose first n iterates are disjoint (a small compact neighborhood of any point will do). By minimality, K is n -good and N -spanning for some N . Let us show that if $N \geq (d+2)n$, then there exists an enlargement of K that is n -good and $N-1$ -spanning.

Assume K is not $N-1$ -spanning. Let K' be a compact neighborhood of K that is n -good, N -spanning and d -mild, and also not $N-1$ -spanning. Let K_* , resp. K'_* , be the set of x in K , resp. K' , such that for each i with $1 \leq i < N$, the point $f^i(x)$ does not belong to K , resp. K' . Let Y be the non-empty set $\overline{K'_*}$. Notice that Y is contained in K_* and thus is N -good. To simplify notation, replace K with K' .

If $x \in Y$ and $f^j(x) \in K$ with $1 \leq j \leq N-1$, then $f^j(x) \in \partial K$. For $x \in Y$, let $J(x)$ be the set of all j with $1 \leq j \leq N-1$ and $f^j(x) \notin K$. Let $I(x)$ be the set of all i with $1 \leq i \leq d+1$ such that $\{in, in+1, \dots, (i+1)n-1\} \subset J(x)$. Since K is d -mild, $\{1, \dots, N-1\} \setminus J(x)$ has at most d elements and $I(x) \neq \emptyset$ for every $x \in Y$.

If for each $x \in Y$ we choose a non-empty subset $\tilde{I}(x)$ of $I(x)$, then the formula

$$K' = K \cup \bigcup_{x \in Y} \bigcup_{i \in I_\varepsilon(x)} f^{in}(x)$$

defines an n -good and $N - 1$ -spanning set. For K' to be compact, we need to guarantee that

$$(9) \quad \text{for any } i, \text{ the set } \{x \in Y : \tilde{I}(x) \ni i\} \text{ is closed.}$$

For $x \in Y$ and $\varepsilon > 0$, let $B(x, \varepsilon) = \{y \in Y : d(y, x) < \varepsilon\}$. Let $J_\varepsilon(x) = \bigcap_{y \in B(x, \varepsilon)} J(y)$ and $I_\varepsilon(x) = \bigcap_{y \in B(x, \varepsilon)} I(y)$. For any $x \in Y$, there is $\varepsilon(x) > 0$ such that $J_{\varepsilon(x)}(x) = J(x)$, and hence $I_{\varepsilon(x)}(x) = I(x)$. Since Y is compact, it can be covered by finitely many balls $B(x_k, \varepsilon(x_k))$. Let ε be a Lebesgue number for this cover. Then, for each $x \in Y$, there exists some k such that $B(x, \varepsilon) \subset B(x_k, \varepsilon(x_k))$ and therefore

$$I_\varepsilon(x) = \bigcap_{y \in B(x, \varepsilon)} I(y) \supset \bigcap_{y \in B(x_k, \varepsilon(x_k))} I(y) = I_{\varepsilon(x_k)}(x_k) = I(x_k) \neq \emptyset.$$

Defining $\tilde{I}(x) = I_\varepsilon(x)$, it is easy to see that (9) is satisfied. Thus there exist an n -good and $N - 1$ -spanning compact set, as we wanted to show. \square

Remark 5.5. Motivated by this result, we pose the following question: If f is a homeomorphism of a compact space (possibly of infinite dimension) without periodic points, is it true that for every $n \in \mathbb{N}$, there exists a compact set whose first n iterates are disjoint and such that finitely many iterates cover the whole space? A related result is Theorem 3.1 in [BC].

Proposition 5.1 follows readily from the previous lemmas.

6. FROM SLOW GROWTH TO INVARIANT SECTIONS FOR THE DISK ACTION

6.1. The Disk Action and the Adjustment Lemma. The group $\text{SL}(2, \mathbb{R})$ acts on the upper half-plane $\mathbb{H} = \{w \in \mathbb{C} : \text{Im } w > 0\}$ as follows:

$$A = \begin{pmatrix} a & b \\ c & d \end{pmatrix} \in \text{SL}(2, \mathbb{R}) \Rightarrow A \cdot w = \frac{aw + b}{cw + d}.$$

(In fact, the action factors through $\text{PSL}(2, \mathbb{R})$.) We fix the following conformal equivalence between \mathbb{H} and the disk $\mathbb{D} = \{z \in \mathbb{C} : |z| < 1\}$:

$$w = \frac{-iz - i}{z - 1} \in \mathbb{H} \leftrightarrow z = \frac{w - i}{w + i} \in \mathbb{D}.$$

Conjugating through this equivalence, we get an action of $\text{SL}(2, \mathbb{R})$ on \mathbb{D} , which we also denote as $(A, z) \mapsto A \cdot z$.

The disk is endowed with the Riemannian metric

$$v \in T_z \mathbb{D} \Rightarrow \|v\|_z = \frac{2|v|}{1 - |z|^2}.$$

Let $d(\cdot, \cdot)$ denote the induced distance function. The group $\text{SL}(2, \mathbb{R})$ acts on \mathbb{D} by isometries.

It can be shown that

$$(10) \quad \|A\| = e^{d(A \cdot 0, 0)/2} \quad \text{for all } A \in \text{SL}(2, \mathbb{R}).$$

In particular, $A \in \text{SO}(2, \mathbb{R})$ iff $A \cdot 0 = 0$.

Let us recall Lemmas 5 and 6 from [ABD]:

Lemma 6.1. *There exists a continuous map $\Phi : \mathbb{D} \times \mathbb{D} \rightarrow \text{SL}(2, \mathbb{R})$ such that $\Phi(p_1, p_2) \cdot p_1 = p_2$ and $\|\Phi(p_1, p_2) - \text{Id}\| \leq e^{d(p_1, p_2)} - 1$.*

Lemma 6.2 (Disk Adjustment). *For every $n \geq 1$, there exists a continuous map $\Psi_n : \mathrm{SL}(2, \mathbb{R})^n \times \mathbb{D}^2 \rightarrow \mathrm{SL}(2, \mathbb{R})^n$ such that if $\Psi_n(A_1, \dots, A_n, p, q) = (\tilde{A}_1, \dots, \tilde{A}_n)$, then*

- a) $\tilde{A}_n \cdots \tilde{A}_1 \cdot p = q$ and
- b) $\|\tilde{A}_i A_i^{-1} - \mathrm{Id}\| \leq \exp\left(\frac{1}{2n} \mathrm{d}(A_n \cdots A_1 \cdot p, q)\right) - 1$ for $1 \leq i \leq n$.

6.2. Construction of an Invariant Section. Let $f : X \rightarrow X$ be a homeomorphism of a compact metric space. If $A \in C(X, \mathrm{SL}(2, \mathbb{R}))$, the cocycle (f, A) , as defined in the introduction, is a skew-product on $X \times \mathrm{SL}(2, \mathbb{R})$. The disk action of $\mathrm{SL}(2, \mathbb{R})$ allows us to consider also another skew-product denoted by $F_{f,A}^{\mathbb{D}} : X \times \mathbb{D} \rightarrow X \times \mathbb{D}$ and given by $F_{f,A}^{\mathbb{D}}(x, z) = (x, A(x) \cdot z)$. The existence of an invariant section for $F_{f,A}^{\mathbb{D}}$ is easily seen to be equivalent to (f, A) being conjugate to a cocycle of rotations. Indeed, if Φ is given by Lemma 6.1 and $B(x) = \Phi(z(x), 0)$, then $B(f(x))A(x)B(x)^{-1}$ is a rotation for every x .

If $A \in C(X, \mathrm{SL}(2, \mathbb{R}))$ (or, more, generally, $A(x)$ is a 2×2 real matrix depending continuously on $x \in X$), we denote $\|A\|_{\infty} = \sup_{x \in X} \|A(x)\|$.

Proposition 6.3. *Let $f : X \rightarrow X$ be a homeomorphism of a compact metric space that admits a minimal non-periodic finite-dimensional factor, and let $\Lambda \subset X$ be a (possibly empty) compact invariant set. For every $C, \varepsilon > 0$, there exists $\gamma > 0$ such that the following holds. Suppose that $A : X \rightarrow \mathrm{SL}(2, \mathbb{R})$ is continuous,*

$$\|A\|_{\infty} < C \quad \text{and} \quad \lim_{n \rightarrow \infty} \sup_{x \in X} \frac{1}{n} \log \|A^n(x)\| < \gamma,$$

and that $F_{f|\Lambda, A|\Lambda}^{\mathbb{D}}$ admits an invariant section z . Then $A|\Lambda$ extends to a continuous map $\tilde{A} : X \rightarrow \mathrm{SL}(2, \mathbb{R})$ such that $\|\tilde{A} - A\|_{\infty} < \varepsilon$ and z extends to an invariant section \tilde{z} for $F_{f, \tilde{A}}^{\mathbb{D}}$.

Proof. Let d be given by Proposition 5.1. Let $\gamma > 0$ be such that $e^{[1+2d(d+1)]\gamma} < 1 + \varepsilon/C$. Let n_0 be such that $n \geq n_0$ implies $\|A^n(x)\| \leq e^{n\gamma}$ for every $x \in X$.

Let $M : \mathbb{D} \rightarrow [0, \infty)$ be the hyperbolic distance to 0. Fix some extension of z to an element of $C(X, \mathbb{D})$ and set $C_1 = \sup_{x \in X} M(z(x))$. Choose $n \in \mathbb{N}$ so that $n > \max\{n_0, \gamma^{-1} \log C, C_1/\gamma\}$. Let K be the corresponding d -mild, n -good, N -spanning set given by Proposition 5.1, where $N = (d+2)n - 1$.

We have

$$\begin{aligned} \|(F_{f,A}^{\mathbb{D}})^j\|_M &= \sup_{x \in X} \sup_{w \in \mathbb{D}} |\mathrm{d}(A^j(x) \cdot w, 0) - \mathrm{d}(w, 0)| \quad (\text{by definition}) \\ &\leq \sup_{x \in X} |\mathrm{d}(A^j(x) \cdot 0, 0)| \quad (\text{since } \mathrm{d} \text{ is } \mathrm{SL}(2, \mathbb{R})\text{-invariant}) \\ &\leq 2 \sup_{x \in X} \log \|A^j(x)\| \quad (\text{by (10)}) \\ &\leq 2j\gamma, \quad \text{provided } j \geq n > n_0. \end{aligned}$$

Applying Lemma 4.1 we obtain an almost invariant section \hat{z} for $F_{f,A}^{\mathbb{D}}$, which is supported in $\mathrm{int} K$ and satisfies

$$\sup_{x \in X} M(\hat{z}(x)) \leq \sup_{x \in X} M(z(x)) + d \max_{n \leq j \leq N} \|(F_{f,A}^{\mathbb{D}})^j\|_M.$$

Thus we obtain $M(\hat{z}(x)) \leq C_1 + 2dN\gamma$ for every $x \in X$.

Let Ψ_n be given by Lemma 6.2 and put

$$\begin{aligned} (\tilde{A}(x), \tilde{A}(f(x)), \dots, \tilde{A}(f^{n-1}(x))) &= \\ &= \Psi_n(A(x), A(f(x)), \dots, A(f^{n-1}(x)), \hat{z}(x), \hat{z}(f^n(x))), \end{aligned}$$

for each $x \in K$. This defines \tilde{A} on $\bigsqcup_{m=0}^{n-1} f^m(K)$. We let $\tilde{A} = A$ on the rest of X . Clearly \tilde{A} is a continuous extension of $A|\Lambda$.

For each $x \in K$, and $1 \leq m \leq n-1$, let $\tilde{z}(f^m(x)) = \tilde{A}^m(x) \cdot \hat{z}(x)$. This defines \tilde{z} on $\bigsqcup_{m=1}^{n-1} f^m(K)$. We let $\tilde{z} = \hat{z}$ on the rest of X . It is easy to see that \tilde{z} is a continuous extension of $z|\Lambda$ and that \tilde{z} is an invariant section for \tilde{A} .

To complete the proof, we need to check that \tilde{A} and A are $\|\cdot\|_\infty$ -close. We have

$$\begin{aligned} d(A^n(x) \cdot \hat{z}(x), \hat{z}(f^n x)) &\leq d(A^n(x) \cdot \hat{z}(x), 0) + d(0, \hat{z}(f^n x)) \\ &\leq d(\hat{z}(x), 0) + 2 \log \|A^n(x)\| + d(0, \hat{z}(f^n x)) \\ &\leq 2(C_1 + 2dN\gamma) + 2n\gamma \\ &\leq [4 + 2d(d+2)]n\gamma. \end{aligned}$$

Hence, using Lemma 6.2,

$$\begin{aligned} \|\tilde{A} - A\|_\infty &\leq C\|\tilde{A}A^{-1} - \text{Id}\|_\infty \\ &\leq C \sup_{x \in X} [\exp [\frac{1}{2n}d(A^n(x) \cdot \hat{z}(x), \hat{z}(f^n x))] - 1] \\ &\leq C(e^{[2+d(d+2)]\gamma} - 1) \\ &< \varepsilon. \end{aligned}$$

□

Let us obtain Theorem 5 from the proposition just proved:

Proof of Theorem 5. Assume that $f : X \rightarrow X$ is a minimal homeomorphism with a finite-dimensional non-periodic factor. We claim that if X is finite dimensional, or there exist at most countably many ergodic invariant measures for f , then a generic $A \in C(X, \text{SL}(2, \mathbb{R})) \setminus \mathcal{UH}$ satisfies

$$(11) \quad \lim_{n \rightarrow \infty} \sup_{x \in X} \frac{1}{n} \log \|A^n(x)\| = 0.$$

In the first case, this is the main result of [AB]. In the second case, [Bo] shows that for each ergodic invariant measure μ , for a generic $A \in C(X, \text{SL}(2, \mathbb{R}))$, the Lyapunov exponent relative to μ vanishes, so taking the (at most countable) intersection, we conclude that for a generic $A \in C(X, \text{SL}(2, \mathbb{R}))$, the Lyapunov exponent vanishes simultaneously for all ergodic invariant measures, and it follows from Proposition 1 in [AB] that (11) holds.

Applying Proposition 6.3 with $\Lambda = \emptyset$, we conclude that any cocycle $A \in C(X, \text{SL}(2, \mathbb{R})) \setminus \mathcal{UH}$ can be approximated by some \tilde{A} that is conjugate to a cocycle of rotations. □

7. SOLVING THE COHOMOLOGICAL EQUATION

From now on, we assume as usual that $f : X \rightarrow X$ is strictly ergodic with invariant probability measure μ and has a non-periodic finite-dimensional factor.

Lemma 7.1. *Let T be a locally compact separable metric space and let $T^* \subset T$ be a closed subset (possibly empty). Let $\varphi \in C^0(T \times X, \mathbb{R})$ and $\psi \in C(T^* \times X, \mathbb{R})$ be such that $\varphi(t, x) = \psi(t, f(x)) - \psi(t, x) + c(t)$ for $(t, x) \in T^* \times X$, where $c(t) = \int \varphi(t, \cdot) d\mu$. Then, for every continuous function $\varepsilon : T \rightarrow \mathbb{R}_+$, there exist $\tilde{\varphi} \in C(T \times X, \mathbb{R})$, coinciding with φ on $T^* \times X$, and $\tilde{\psi} \in C(T \times X, \mathbb{R})$, coinciding with ψ on $T^* \times X$, such that $|\tilde{\varphi}(t, x) - \varphi(t, x)| < \varepsilon(t)$ and $\tilde{\varphi}(t, x) = \tilde{\psi}(t, f(x)) - \tilde{\psi}(t, x) + c(t)$ for all $(t, x) \in T \times X$.*

The proof of this proposition is somewhat similar to Proposition 6.3, using the abstract Lemma 4.1 and Proposition 5.1.

Proof. Let us assume first that T is compact. In this case, we may replace the function ε by a positive number. Let the functions φ and ψ be given. Let $S_j(t, x) = \sum_{m=0}^{j-1} \varphi(t, f^m(x))$ and $c(t) = \int \varphi(t, x) d\mu(x)$. Using that T is compact, it is easy to see that $\frac{1}{j} S_j(t, x)$ converges to $c(t)$ uniformly over $(t, x) \in T \times X$.

Let d be given by Proposition 5.1. Choose $\delta > 0$ with $[2d(d+2)+1]\delta < \varepsilon$, and choose j_0 such that $\left| \frac{1}{j} S_j(t, x) - c(t) \right| < \delta$, uniformly for $(t, x) \in T \times X$, and $j \geq j_0$. Then choose $n \geq j_0$ such that $\frac{2\|\psi\|_\infty}{n} + [2d(d+2)+1]\delta < \varepsilon$. (Here $\|\psi\|_\infty$ is meant to be 0 in the case that $T^* = \emptyset$.) By Proposition 5.1, there is a compact set K that is d -mild, n -good, and N -spanning, where $N = (d+2)n - 1$.

Let $\hat{X} = T \times X$, and $\hat{f} : \hat{X} \ni (t, x) \mapsto (t, f(x))$. Consider the skew-product on $\hat{X} \times \mathbb{R}$ over \hat{f} given by

$$F = F_{\hat{f}, \varphi}^{\mathbb{R}} : (t, x, w) \mapsto (t, f(x), w + \varphi(t, x) - c(t)).$$

Let $\Lambda = T^* \times X$. By assumption, the map $\psi \in C(\Lambda, \mathbb{R})$ is an invariant section of the skew product $F|_{\Lambda \times \mathbb{R}}$ over $\hat{f}|_{\Lambda}$. Let $\psi_0 \in C(\hat{X}, \mathbb{R})$ be a continuous extension of ψ , so that $\|\psi_0\|_\infty = \|\psi\|_\infty$. (If $T^* = \emptyset$, take $\psi_0 \equiv 0$.) Then ψ_0 is an almost invariant section for the skew-product F supported outside Λ . Applying Lemma 4.1 with the pair $(Y, M) = (\mathbb{R}, |\cdot|)$, we obtain an almost invariant section $\hat{\psi} \in C(\hat{X}, \mathbb{R})$ supported in $\text{int } K$ (i.e., $\varphi(t, x) = \hat{\psi}(t, f(x)) - \hat{\psi}(t, x) + c(t)$, provided $x \notin \text{int } K$) that coincides with ψ_0 (and hence with ψ) on Λ , and

$$\begin{aligned} \|\hat{\psi}\|_\infty &\leq \|\psi_0\|_\infty + d \max_{j \in [n, N]} \|F^j\|_M \\ &= \|\psi\|_\infty + d \max_{j \in [n, N]} \sup_{(t, x)} |S_j(t, x) - jc(t)| \\ &< \|\psi\|_\infty + dN\delta. \end{aligned}$$

Let us define a function $\tilde{\varphi} : \hat{X} \rightarrow \mathbb{R}$. For each $t \in T$, $x \in K$, and $0 \leq m < n$, set

$$\tilde{\varphi}(t, f^m(x)) = \varphi(t, f^m(x)) + \frac{\hat{\psi}(t, f^n(x)) - \hat{\psi}(t, x) - S_n(t, x) + nc(t)}{n}.$$

This defines $\tilde{\varphi}$ on $T \times \bigsqcup_{m=0}^{n-1} f^m(K)$; let it equal φ in the complement of this set. Then $\tilde{\varphi}$ is continuous and

$$\|\tilde{\varphi} - \varphi\|_\infty \leq \frac{2\|\psi\|_\infty}{n} + [2d(d+2)+1]\delta < \varepsilon.$$

Notice that if $t \in T^*$, then $\hat{\psi}(t, x) = \psi_0(t, x) = \psi(t, x)$ and thus $\tilde{\varphi}(t, x) = \varphi(t, x)$.

Let us define another function $\tilde{\psi} : \hat{X} \rightarrow \mathbb{R}$. For each $t \in T$, $x \in K$, and $0 \leq m < n$, set

$$\tilde{\psi}(t, f^m(x)) = \hat{\psi}(t, x) + \sum_{j=0}^{m-1} [\tilde{\varphi}(t, f^j(x)) - c(t)].$$

This defines $\tilde{\psi}$ on $T \times \bigsqcup_{m=0}^{n-1} f^m(K)$; let it equal $\hat{\psi}$ in the complement of this set. Then $\tilde{\psi}$ is continuous and $\tilde{\varphi}(t, x) = \tilde{\psi}(t, f(x)) - \tilde{\psi}(t, x) + c(t)$ for all $(t, x) \in T \times X$. Also, $\tilde{\psi}(t, x) = \psi(t, x)$ when $t \in T^*$. So the functions $\tilde{\varphi}$ and $\tilde{\psi}$ have all the desired properties.

This concludes the proof in the case where T is compact. Let us now consider the general case. The hypotheses on T imply that there exists an exhaustion of T in the sense that $T = \bigcup_{i=1}^{\infty} T_i$ with T_i compact and $T_i \subset \text{int } T_{i+1}$.

We initially define $\tilde{\varphi}$ and $\tilde{\psi}$ on $T_1 \times X$, by applying the compact case with the data $T_{\text{new}} = T_1$, $T_{\text{new}}^* = T_1 \cap T^*$, $\varphi_{\text{new}} = \varphi|_{T_1}$, $\psi_{\text{new}} = \psi|(T_1 \cap T^*)$, $\varepsilon_{\text{new}} = \varepsilon|_{T_1}$: this yields functions $\tilde{\varphi}_{\text{new}}, \tilde{\psi}_{\text{new}} : T_1 \times X \rightarrow \mathbb{R}$, which we take as the definition of $\tilde{\varphi}$

and $\tilde{\psi}$ on $T_1 \times X$. By construction, all requirements are satisfied when we restrict considerations to $T_1 \times X$.

Assume we have already defined $\tilde{\varphi}$ and $\tilde{\psi}$ on $T_i \times X$ so that all requirements are satisfied when considerations are restricted to $T_i \times X$. Choose a continuous function $\varphi_i : T_{i+1} \times X \rightarrow \mathbb{R}$ that coincides with $\tilde{\varphi}$ on $T_i \times X$ and with φ on $(T^* \cap T_{i+1}) \times X$, such that $\int \varphi_i(t, x) d\mu(x) = c(t)$ for every $t \in T_{i+1}$ and satisfying $|\varphi_i(t, x) - \varphi(t, x)| < \varepsilon(t)$ for every $(t, x) \in T_{i+1} \times X$. Now apply the compact case again with data $T_{\text{new}} = T_{i+1}$, $T_{\text{new}}^* = T_i \cup (T_{i+1} \cap T^*)$, $\varphi_{\text{new}} = \varphi_i$, $\psi_{\text{new}} = \tilde{\psi}|(T_i \cup (T_{i+1} \cap T^*))$ and $\varepsilon_{\text{new}}(t) = \varepsilon(t) - \sup_{x \in X} |\varphi_i(t, x) - \varphi(t, x)|$. This yields functions $\tilde{\varphi}_{\text{new}}, \tilde{\psi}_{\text{new}} : T_{i+1} \times X \rightarrow \mathbb{R}$, which we take as the definitions of $\tilde{\varphi}$ and $\tilde{\psi}$ on $T_{i+1} \times X$. Notice that, by construction, the new definitions extend the previous ones, and satisfy all the requirements when we restrict considerations to $T_{i+1} \times X$.

By induction, this procedure yields functions $\tilde{\varphi}, \tilde{\psi} : T \times X \rightarrow \mathbb{R}$. These functions are continuous since they locally coincide with a continuous function obtained in a finite stage of the induction (here we use that $\bigcup \text{int } T_i = T$), and satisfy all the other requirements over $T \times X$ (since those can be also verified at a finite stage of the induction). \square

Now we are able to give the:

Proof of Theorem 8. Let $A \in \mathcal{UH}$, and let $D \in \text{SL}(2, \mathbb{R})$ be such that $L(D) = L(A)$. Let us indicate

$$\Delta(r) = \begin{pmatrix} e^r & 0 \\ 0 & e^{-r} \end{pmatrix}.$$

We can assume that $D = \Delta(c)$ for some $c \in \mathbb{R}$. It is not difficult (see, e.g., the proof of Lemma 4 in [ABD]) to show that there are continuous maps $B : X \rightarrow \text{PSL}(2, \mathbb{R})$ and $\varphi : X \rightarrow \mathbb{R}$ such that $A(x) = B(fx)^{-1} \Delta(\varphi(x)) B(x)$ (in $\text{PSL}(2, \mathbb{R})$) and $\int \varphi d\mu = c$.

Apply Lemma 7.1 (with $T = (0, 1]$ and $T^* = \emptyset$) to find continuous functions $\tilde{\varphi}, \tilde{\psi} : (0, 1] \times X \rightarrow \mathbb{R}$ such that

$$\begin{aligned} |\tilde{\varphi}(t, x) - \varphi(x)| &< \varepsilon(t) \quad \text{where } \varepsilon(t) \rightarrow 0 \text{ as } t \rightarrow 0, \\ \tilde{\varphi}(t, x) &= \tilde{\psi}(t, fx) - \tilde{\psi}(t, x) + c. \end{aligned}$$

Define

$$A_t(x) = B(fx)^{-1} \Delta(\varphi(t, x)) B(x) \quad \text{and} \quad B_t(x) = \Delta(-\tilde{\psi}(t, x)) B(x).$$

Then $A_t \rightarrow A$ as $t \rightarrow 0$ and $A_t(x) = B_t(fx) D B_t(x)^{-1}$, as desired. \square

8. EXTENDING FAMILIES OF COCYCLES: SMALL LYAPUNOV EXPONENTS

Let T be a topological space. We say that A_t , $t \in T$ is a *continuous family of cocycles* if $A_t \in C(X, \text{SL}(2, \mathbb{R}))$ for every $t \in T$ and the map $T \ni t \mapsto A_t \in C(X, \text{SL}(2, \mathbb{R}))$ is continuous. If \mathcal{B} is a subset of $C(X, \text{SL}(2, \mathbb{R}))$ that contains $\{A_t : t \in T\}$ and such that the map $T \ni t \mapsto A_t \in \mathcal{B}$ is homotopic to a constant, then we say that that *the family of cocycles A_t is contractible in \mathcal{B}* . For example, a continuous family of cocycles A_t , $t \in \partial[0, 1]^p$ taking values in a set \mathcal{B} is contractible in \mathcal{B} if and only if there is an extended continuous family of cocycles A_t , $t \in [0, 1]^p$ also taking values in \mathcal{B} .

Given a continuous family of cocycles A_t , $t \in T$, we say that z_t , $t \in T$ is a *continuous family of sections* if each $t \in T$ corresponds to an invariant section $z_t \in C(X, \mathbb{D})$ for the skew-product $F_{f, A_t}^{\mathbb{D}}$ (see §6.2), and the map $T \ni t \mapsto z_t \in C(X, \mathbb{D})$ is continuous.

The parameter spaces that we use next are the p -dimensional cube $[0, 1]^p$ and its boundary $\partial[0, 1]^p$. The goal of this section is to prove the following result:

Lemma 8.1. *Let $p \geq 1$ and let $\mathcal{B} \subset C(X, \mathrm{SL}(2, \mathbb{R}))$ be a bounded open set. Let A_t , $t \in \partial[0, 1]^p$ be a continuous family of cocycles that is contractible in \mathcal{B} , and such that $L(A_t) = 0$ for all $t \in \partial[0, 1]^p$. Then, for any $\gamma > 0$, we can find a continuous family of cocycles \hat{A}_t , $t \in [0, 1]^p$ taking values in \mathcal{B} and extending the family A_t , $t \in \partial[0, 1]^p$, such that $L(\hat{A}_t) < \gamma$ for every $t \in [0, 1]^p$.*

The proof of Lemma 8.1 requires several preliminaries. A measurable set $Y \subset X$ is called a *coboundary* if there is a measurable set Z such that $Y = Z \Delta f^{-1}(Z)$ μ -mod 0. Non-coboundaries are relatively easy to find:

Lemma 8.2 (Corollary 3.5 from [K]). *If $Z \subset X$ is a positive measure set, then there is a measurable set $Y \subset Z$ that is not a coboundary.*

Our interest in non-coboundaries is due to the following fact, already used in [K]:

Lemma 8.3. *Let $A : X \rightarrow \mathrm{SL}(2, \mathbb{R})$ be such that $\log \|A\|$ is μ -integrable. Assume that there exist two measurable maps $e^1, e^2 : X \rightarrow \mathbb{P}^1$ such that*

$$e^1(x) \neq e^2(x), \quad A(x) \cdot \{e^1(x), e^2(x)\} = \{e^1(fx), e^2(fx)\} \quad \text{for } \mu\text{-almost every } x.$$

If the set $Y = \{x \in X : A(x) \cdot e^1(x) = e^2(fx)\}$ is not a coboundary, then the Lyapunov exponent $L(A)$ vanishes.

Proof. Assume that $L(A) > 0$, and consider the two Oseledets directions $e^+, e^- : X \rightarrow \mathbb{P}^1$. It follows from Oseledets's Theorem that $\{e^1(x), e^2(x)\} = \{e^+(x), e^-(x)\}$ for almost every x . Let W be the set of x where $e^1(x) = e^+(x)$. Then $W \Delta f^{-1}(W)$ is precisely the set Y . Hence it is a coboundary. \square

If $Z \subset X$ is a measurable set of positive measure, the *first-return map* $f_Z : Z \rightarrow Z$ is defined as $f_Z(x) = f^{r_Z(x)}(x)$, where $r_Z(x)$ is the least positive integer n such that $f^n(x) \in Z$. The probability measure $\frac{\mu}{\mu(Z)}$ is invariant and ergodic with respect to f_Z . If $A : X \rightarrow \mathrm{SL}(2, \mathbb{R})$ is such that $\log \|A\| \in L^1(\mu)$, then we define the *induced cocycle* $A_Z : Z \rightarrow \mathrm{SL}(2, \mathbb{R})$ by $A_Z(x) = A^{r_Z(x)}(x)$.

Lemma 8.4 (cf. Lemma 2.2 from [K]). *The Lyapunov exponent of the cocycle A_Z over the dynamics f_Z and with respect to the ergodic measure $\frac{\mu}{\mu(Z)}$ is equal to $L(A)/\mu(Z)$.*

A well-known theorem by Lusin states that a measurable function can be altered in a set of arbitrarily small measure so that it becomes continuous. We will need a parameterized version of this:

Lemma 8.5. *Assume T is a compact metric space. Let $\psi : T \times X \rightarrow \mathbb{R}$ be a Borel measurable function such that for every $t \in T$, there is a full measure subset of X consisting of points x with the property that the function $t' \mapsto \psi(t', x)$ is continuous on a neighborhood of t . Then for every $\beta > 0$, there exists a continuous function $\phi : T \times X \rightarrow \mathbb{R}$ such that for every $t \in T$, the points $x \in X$ where $\phi(t, x) \neq \psi(t, x)$ form a set of measure less than β .*

Proof. For each $t \in T$ and $r > 0$, let $\bar{B}(t, r)$ be the closed ball of radius r centered at t , and

$$G(t, r) = \{x \in X : t' \mapsto \psi(t', x) \text{ is continuous on } \bar{B}(t, r)\}.$$

For every $t \in T$, there is $\rho(t) > 0$ such that $\mu(G(t, \rho(t))) > 1 - \beta/3$. Take a cover of T by finitely many balls $B_i = \bar{B}(t_i, \rho(t_i))$. Now, for each i , let $K_i \subset G(t_i, \rho(t_i))$ be compact with $\mu(K_i) > 1 - 2\beta/3$. Consider the mapping $\Psi_i : K_i \rightarrow C(B_i, \mathbb{R})$ given by $t' \mapsto (x \in B_i \mapsto \psi(t', x))$. Take a subset $K'_i \subset K_i$ with $\mu(K'_i) > 1 - \beta$ such that $\Psi_i|K'_i$ is continuous. Let $L = \bigcup_i B_i \times K'_i$. Then L is compact, ψ restricted to L is continuous, and for each t , the set of x such that (t, x) belongs to L has measure at

least $1 - \beta$. Let $\phi : T \times X \rightarrow \mathbb{R}$ be a continuous extension of $\psi|L$, given by Tietze's Extension Theorem. Then ϕ has the required properties. \square

We say that A_t , $t \in T$ is an L^∞ family of cocycles if $A_t \in L^\infty(X, \text{SL}(2, \mathbb{R}))$ for every $t \in T$, the map $(t, x) \mapsto A_t(x)$ is measurable and $\sup_t \|A_t\|_\infty$ is finite.

We say that A_t , $t \in T$ is an almost-continuous family of cocycles if it is an L^∞ family, T is a compact metric space, and for every $t \in T$, there is a full measure set $G_t \subset X$ such that for each $x \in G_t$, the mapping $t' \mapsto A_{t'}(x)$ is continuous in a neighborhood of t .

We define a weak metric on the set $L^\infty(X, \text{SL}(2, \mathbb{R}))$ as follows:

$$d_w(A, B) = \inf \{ \beta > 0 : \mu\{x \in X : \|A(x) - B(x)\| > \beta\} < \beta \}.$$

Lemma 8.6. *Assume A_t , $t \in T$ is an almost-continuous family of cocycles. Let $\lambda = \sup_{t \in T} L(A_t)$, and fix $M > \sup_t \|A_t\|_\infty$. Then:*

- a) $\lim_{n \rightarrow \infty} \sup_{t \in T} \frac{1}{n} \int_X \log \|A_t^n\| d\mu = \lambda$.
- b) *For every $\gamma > \lambda$, there exists $\beta > 0$ with the following properties. Assume that B_t , $t \in T$ is an L^∞ family of cocycles such that $\|B_t\|_\infty < M$ and $d_w(B_t, A_t) < \beta$ for every $t \in T$. Then $L(B_t) < \gamma$ for every t .*

Proof. The limit in part (a) exists by subadditivity, and is obviously bigger or equal than λ . Let us prove the reverse inequality.

Fix $\eta > 0$ and take $t \in T$. Let $m \in \mathbb{N}$ be such that $\frac{1}{m} \int \log \|A_t^m\| d\mu < \lambda + \eta$. Since A_t is an almost-continuous family of cocycles, so is A_t^m . Hence for any $\varepsilon > 0$, there exists $\delta > 0$ such that

$$\mu\{x \in X : \|A_{t'}^m(x) - A_t^m(x)\| < \varepsilon \text{ for all } t' \in B(t, \delta)\} > 1 - \varepsilon.$$

Since the family is essentially bounded, we can find $\delta = \delta(t) > 0$ such that $\frac{1}{m} \int \log \|A_{t'}^m\| d\mu < \lambda + 2\eta$ for every $t' \in B(t, \delta)$. It follows from subadditivity that there exists $\bar{m} = \bar{m}(t) \in \mathbb{N}$ such that

$$\frac{1}{n} \int \log \|A_{t'}^n\| d\mu < \lambda + 3\eta \quad \text{for all } t' \in B(t, \delta) \text{ and } n \geq \bar{m}.$$

Now take a cover of T by finitely many balls $B(t_i, \delta(t_i))$. Let $n_0 = \max_i \bar{m}(t_i)$. Then

$$(12) \quad \frac{1}{n} \int \log \|A_t^n\| d\mu < \lambda + 3\eta \quad \text{for every } t \in T \text{ and } n \geq n_0.$$

Since $\eta > 0$ is arbitrary, part (a) of the Lemma follows.

Let $\varepsilon > 0$ be small. There exists $\beta > 0$ such that if B_t , $t \in T$ is an L^∞ family of cocycles with $d_w(B_t, A_t) < \beta$ for each t , then $d_w(B_t^{n_0}, A_t^{n_0}) < \varepsilon$. Now, if in addition $\|B_t\|_\infty < M$ and ε was chosen small enough, it follows from (12) that

$$\frac{1}{n_0} \int \log \|B_t^{n_0}\| d\mu < \lambda + 4\eta,$$

which implies $L(B_t) < \lambda + 4\eta$ for every $t \in T$. This proves part (b). \square

After these preliminaries, we can give the:

Proof of Lemma 8.1. Let us first explain informally the steps of the proof:

- a) We select a bounded cocycle D in \mathcal{B} and a small tower for f along which some products of D are rotations.
- b) By assumption, the given family A_t of cocycles can be extended to a family defined in the whole cube, which we also denote by A_t .

- c) In the crucial step, we find an almost-continuous family of cocycles B_t that has zero exponents for all values of the parameter and is uniformly very close to A_t everywhere except in the tower. This is done using Lemma 8.3.
- d) We use Lemma 8.5 to approximate B_t by a continuous family of cocycles \tilde{A}_t ; by Lemma 8.6 we get $L(\tilde{A}_t)$ small for every t .
- e) Since the original family of cocycles A_t , $t \in \partial[0, 1]^p$ has zero exponents, we can interpolate it with the \tilde{A}_t family near the boundary of the cube and obtain a continuous family \hat{A}_t , $t \in [0, 1]^p$ that coincides with the initial family in the boundary, and (by semicontinuity) has small exponent in the whole cube.

Now let us give the actual proof.

Since the set \mathcal{B} is bounded and $L(A_t) = 0$ for every $t \in \partial[0, 1]^p$, by part (b) of Lemma 8.6 there is $\beta > 0$ such that whenever A'_t , $t \in \partial[0, 1]^p$ is an L^∞ family of cocycles taking values in \mathcal{B} and with $d_w(A'_t, A_t) < \beta$, we have $L(A'_t) < \gamma$ for every $t \in \partial[0, 1]^p$.

Since the family A_t , $t \in \partial[0, 1]^p$ is contractible in \mathcal{B} , we can extend it to a continuous family A_t , $t \in [0, 1]^p$ also taking values in \mathcal{B} .

Let $\varepsilon > 0$ be very small. (Actual smallness requirements will emerge several times along of the proof.)

Let $t_c = (\frac{1}{2}, \dots, \frac{1}{2})$ be the center of the cube. We may assume that $A_{t_c} = A_t$ for some $t \in \partial[0, 1]^p$. Using Theorem 5, choose a cocycle D close to A_{t_c} that has an invariant section $z : X \rightarrow \mathbb{D}$. Making if necessary a small perturbation, we can assume that z assumes the same value z_0 in an open set. By conjugating everything with some fixed matrix (and hence changing \mathcal{B} , γ , β accordingly), we can assume that $z_0 = 0$. Let N be a large integer. Take a non-empty open set $Z \subset X$ such that its iterates in times $0 = n_0 < n_1 < \dots < n_N$ are wholly contained in $\{z = 0\}$, and moreover $Z \cap f^i(Z) = \emptyset$ for $0 < i \leq n_N + 1$. Notice that if $x \in Z$, then $D^{n_i}(x) \in \text{SO}(2, \mathbb{R})$ for $0 \leq i \leq N$. By a further perturbation (and reducing Z if necessary), we can assume that in addition to the previous properties, $D|f^i(Z)$ is constant for $0 \leq i \leq n_N$, and equal to the identity matrix for $i = n_N$. Reducing Z further, we assume that the tower $\bigcup_{i=0}^{n_N-1} f^i(Z)$ has measure less than $\beta/2$. Let $W = f^{n_N}(Z)$ and $F = f^{n_N} : Z \rightarrow W$. For $x \in W$, let $\ell(x) > 0$ be minimal with $f^{\ell(x)}(x) \in Z$, and let $G : W \rightarrow Z$ be given by $G(x) = f^{\ell(x)}(x)$.

Since the set \mathcal{B} is open, we can modify the family A_t , $t \in [0, 1]^p$ near $t = t_c$ so that A_{t_c} equals D .

Let $P_t : W \rightarrow \text{SL}(2, \mathbb{R})$ be given by $P_t(x) = A_t^{\ell(x)}(x)$.

Claim. There exists an L^∞ family \tilde{P}_t , $t \in [0, 1]^p$ with the following properties:

- a) $\|P_t^{-1}\tilde{P}_t - \text{Id}\|_\infty < \varepsilon$ for every t .
- b) $(t, x) \mapsto \tilde{P}_t(x)$ is finite-valued and takes values in $\text{SL}(2, \mathbb{R}) \setminus \text{SO}(2, \mathbb{R})$.
- c) For every t , there is a full measure subset of W consisting of points x such that the map $t' \mapsto \tilde{P}_{t'}(x)$ is constant in a neighborhood of t .

Proof. Since f is minimal and Z is open, $\ell(x)$ is bounded and therefore so is the function $P_t(x)$. Choose then a large finite set of $\text{SL}(2, \mathbb{R}) \setminus \text{SO}(2, \mathbb{R})$ approximating its image.

Since (X, μ) is a non-atomic Lebesgue space, we can assume X is the unit interval and μ is Lebesgue measure. Let \mathcal{Q} be a partition of \mathbb{R}^p into small dyadic cubes. Fix a vector $v_0 \in \mathbb{R}^p$ whose coordinates are independent mod \mathbb{Q} . For each $x \in [0, 1]$, let $\mathcal{Q}(x)$ be the restriction of the partition $xv_0 + \mathcal{Q}$ to the unit cube. Then every $t \in [0, 1]^p$ belongs to the interior (relative to $[0, 1]^p$) of an element of the partition $\mathcal{Q}(x)$ for all but finitely many points $x \in X$. Take $\tilde{P}_t(x)$ to be constant in each element of $\mathcal{Q}(x)$, and taking values in the large finite set. This proves the claim. \square

For each $x \in W$, let $e_t^s(x)$ and $e_t^u(x)$ be respectively the most contracted and the most expanded directions by $\tilde{P}_t(x)$; these two directions are uniquely defined because the matrix is not a rotation, and are orthogonal. Next define $e_t^s(x)$ and $e_t^u(x)$ for $x \in Z = G(W)$ by

$$e_t^*(x) = \tilde{P}_t(G^{-1}(x)) \cdot e_t^*(G^{-1}(x)), \quad * \in \{s, u\}.$$

These two directions are also orthogonal.

There are only finitely many possible values for $e_t^u(x)$ and $e_t^s(x)$; let $\mathcal{F} \subset \mathbb{P}^1$ be the set of all of them. Define a function $\angle : \mathcal{F} \times \mathcal{F} \rightarrow [0, \pi)$ such that $R_{\angle(x,y)} \cdot x = y$.

Apply Lemma 8.2 and choose some set of positive measure $Y \subset Z$ that is not a coboundary. Let

$$\theta(t, x) = \begin{cases} \angle(e_t^u(x), e_t^s(F(x))) & \text{if } x \in Y, \\ \angle(e_t^u(x), e_t^u(F(x))) & \text{if } x \in Z \setminus Y. \end{cases}$$

Notice that for every t and for every x in a full measure subset of X (depending on t), the function $t' \mapsto \theta(t', x)$ is constant in a neighborhood of t .

Define

$$B_t(x) = \begin{cases} D(x) \cdot R_{\theta(t, f^{-n_i}x)/N} & \text{if } x \in f^{n_i}(Z), 0 \leq i < N, \\ D(x) & \text{if } x \in f^j(Z), j \in \{0, 1, \dots, n_N\} \setminus \{n_0, n_1, \dots, n_{N-1}\}, \\ A_t(x) \cdot [P_t(x)]^{-1} \cdot \tilde{P}_t(x) & \text{if } x \in W, \\ A_t(x) & \text{otherwise.} \end{cases}$$

Let us list some properties of the family of cocycles B_t , $t \in [0, 1]^p$. First, it is almost-continuous. Second,

$$(13) \quad L(B_t) = 0 \quad \text{for every } t \in [0, 1]^p.$$

Indeed, the induced cocycle $(B_t)_Z(x)$ is given by

$$C_t(x) = \tilde{P}_t(F(x))R_{\theta(t,x)}, \quad x \in Z.$$

(Of course, this cocycle is regarded over the dynamics $G \circ F$, which is the first return map to Z .) The pair of measurable directions e_t^u , e_t^s is preserved by the action of the cocycle; moreover the directions are swapped precisely for the points in Y . Thus Lemma 8.3 gives that for each t , the measurable cocycle $(G \circ F, C_t)$ has zero Lyapunov exponent (with respect to the ergodic invariant measure $\frac{1}{\mu(Z)}\mu$). Thus (13) follows from Lemma 8.4.

A third property of the family B_t is that for each t , $B_t(x)$ is close to $A_t(x)$ for every x outside a set of small measure. Let us be more precise. Let $\tau(t, x)$ equal t_c if x is in the tower $\bigcup_{i=0}^{n_N} f^i(Z)$ (which has measure less than $\beta/2$), and t otherwise. Then we can write

$$B_t(x) = A_{\tau(t,x)}(x) \cdot E_t(x),$$

where E_t , $t \in [0, 1]^p$ is an almost continuous family of cocycles such that $\|E_t(x) - \text{Id}\| < \varepsilon$ for all t and x . (Here we used that N is large.)

Since $\tau(t, x)$ depends continuously on t , we can apply for instance Lemma 8.5 and find a continuous function $\tilde{\tau} : [0, 1]^p \times X \rightarrow [0, 1]^p$ such that $\mu\{x \in X : \tilde{\tau}(t, x) \neq \tau(t, x)\}$ is very small for every t .

Let \mathbb{B}_ε be the set of $M \in \text{SL}(2, \mathbb{R})$ such that $\|M - \text{Id}\| < \varepsilon$. Since ε is small, this set is homeomorphic to \mathbb{R}^3 . Thus we can apply Lemma 8.5 and find a continuous map $(t, x) \in [0, 1]^p \times X \mapsto \tilde{E}_t(x) \in \mathbb{B}_\varepsilon$ such that $\mu\{x \in X : \tilde{E}_t(x) \neq E_t(x)\}$ is uniformly small for every t .

Now define

$$\tilde{A}_t(x) = A_{\tilde{\tau}(t,x)}(x) \cdot \tilde{E}_t(x).$$

Then \tilde{A}_t is a continuous family of cocycles and (because ε is small) it takes values in \mathcal{B} . We can choose $\tilde{\tau}(t, x)$ and $\tilde{E}_t(x)$ so that $\sup_t d_w(\tilde{A}_t, B_t)$ is as small as desired. Therefore, by (13) and part (b) of Lemma 8.6, we can assume that

$$L(\tilde{A}_t) < \gamma \quad \text{for every } t \in [0, 1]^p.$$

With appropriate choices, we also have that for each $t \in [0, 1]^p$, the set

$$G_t = \{x \in X : \tilde{\tau}(t, x) = t\} \quad \text{has } \mu(G_t) > 1 - \frac{3}{4}\beta.$$

The last step of the proof is to modify the family \tilde{A}_t to make it match with A_t in the boundary of the cube, while keeping the exponent small. Let $\pi : [0, 1]^p \setminus \{t_c\} \rightarrow \partial[0, 1]^p$ be a retraction. Let S be a neighborhood of $\partial[0, 1]^p$ inside $[0, 1]^p$, thin enough so that for every $t \in S$, $\pi(t)$ is close to t and hence $A_{\pi(t)}$ is close to A_t . Let \hat{E}_t , $t \in [0, 1]^p$ be a continuous family taking values in \mathbb{B}_ε , such that $\hat{E}_t(x) = \tilde{E}_t(x)$ for t outside S , and $\hat{E}_t(x) = \text{Id}$ for $t \in \partial[0, 1]^p$. Let $\hat{\tau} : [0, 1]^p \times X \rightarrow [0, 1]^p$ be continuous such that $\hat{\tau}(t, x) = \tilde{\tau}(t, x)$ if $t \notin S$, $\hat{\tau}(t, x) = t$ if $t \in \partial[0, 1]^p$, and moreover $\hat{\tau}(t, x) = t$ for every $t \in [0, 1]^p$ and $x \in G_t$. Define

$$\hat{A}_t(x) = A_{\hat{\tau}(t, x)}(x) \cdot \hat{E}_t(x).$$

Then the family \hat{A}_t , $t \in [0, 1]^p$ takes values in \mathcal{B} and continuously extends A_t , $t \in \partial[0, 1]^p$. For $t \in [0, 1]^p \setminus S$, we have $\hat{A}_t = \tilde{A}_t$ and therefore $L(\hat{A}_t) < \gamma$. For $t \in S$ and $x \in G_t$, we have $\hat{A}_t(x) = A_t(x) \cdot \hat{E}_t(x)$, which is close to $A_{\pi(t)}(x)$. Since ε is small, we have $d_w(\hat{A}_t, A_{\pi(t)}) < \beta$. By definition of β , this implies $L(\hat{A}_t) < \gamma$. This completes the proof of Lemma 8.1. \square

9. EXTENDING FAMILIES OF COCYCLES: BOUNDEDNESS

In this section, we will strengthen Lemma 8.1, obtaining the following two results:

Proposition 9.1. *Let $p \geq 1$, and let $\mathcal{B} \subset C(X, \text{SL}(2, \mathbb{R}))$ be a bounded open set. Let A_t , $t \in \partial[0, 1]^p$ be a continuous family of cocycles that is contractible in \mathcal{B} and admits a continuous family of sections z_t , $t \in \partial[0, 1]^p$. Then z_t can be extended to a continuous family of sections for an extended family of cocycles A_t , $t \in [0, 1]^p$ that also takes values in \mathcal{B} .*

Proposition 9.2. *Let $p \geq 1$, and let $\mathcal{B} \subset C(X, \text{SL}(2, \mathbb{R}))$ be a bounded open set. Let A_t , $t \in \partial[0, 1]^p$ be a continuous family of cocycles that is contractible in \mathcal{B} and such that $A_t \notin \mathcal{UH}$ for every $t \in \partial[0, 1]^p$. Then there exists an extended continuous family of cocycles A_t , $t \in [0, 1]^p$ also taking values in \mathcal{B} and such that the family A_t , $t \in (0, 1)^p$ admits a continuous family of sections.*

Notice that Theorem 4 is an immediate consequence of this proposition with $p = 1$. It also follows that the complement of \mathcal{UH} is locally connected. This shows the statement made in Remark 2.2.

We will first prove Proposition 9.1, and then use it together with Theorem 5 to obtain Proposition 9.2. If $A \in C(X, \text{SL}(2, \mathbb{R}))$ and $\varepsilon > 0$, we write $\mathcal{B}_\varepsilon(A) = \{\tilde{A} : \|\tilde{A} - A\|_\infty < \varepsilon\}$.

Proof of Proposition 9.1. Let A_t , $t \in \partial[0, 1]^p$ be a continuous family of cocycles admitting a continuous family of sections. Assume that this family is contractible in a bounded open set $\mathcal{B} \subset C(X, \text{SL}(2, \mathbb{R}))$, and thus can be extended to a \mathcal{B} -valued family \bar{A}_t , $t \in [0, 1]^p$. Let $\varepsilon > 0$ be such that for each $t \in [0, 1]^p$, the ball with center \bar{A}_t and radius 2ε , indicated by $\mathcal{B}_{2\varepsilon}(\bar{A}_t)$, is contained in \mathcal{B} . Let $\gamma > 0$ be given by Proposition 6.3, where in place of f we take the homeomorphism of $\hat{X} = X \times [0, 1]^p$ given by $\hat{f}(x, t) = (f(x), t)$, and take $\Lambda = X \times \partial[0, 1]^p$, and $C > \sup_{B \in \mathcal{B}} \|B\|_\infty$. Next apply Lemma 8.1 to the family A_t , $t \in \partial[0, 1]^p$ and to the

reduced set $\mathcal{B}' = \bigcup_{t \in [0,1]^p} \mathcal{B}_\varepsilon(\bar{A}_t)$. We obtain a family \hat{A}_t , $t \in [0,1]^p$ with values in \mathcal{B}' such that $L(\hat{A}_t) < \gamma$ for each t . Using compactness of $[0,1]^p$ and subadditivity, it is not hard to see that

$$\lim_{n \rightarrow \infty} \sup_{(x,t) \in X \times [0,1]^p} \frac{1}{n} \log \|\hat{A}_t^n(x)\| < \gamma.$$

Thus we can apply Proposition 6.3 and find a continuous family of cocycles \tilde{A}_t , $t \in [0,1]^p$ that equals \hat{A}_t (and hence A_t) when $t \in \partial[0,1]^p$, such that $\|\tilde{A}_t - \hat{A}_t\|_\infty < \varepsilon$ (and thus $\tilde{A}_t \in \mathcal{B}$) for each t , and that admits a continuous family of sections \tilde{z}_t , $t \in [0,1]^p$ extending the original z_t , $t \in \partial[0,1]^p$. \square

Lemma 9.3. *For every $C > 0$, there is $\bar{\varepsilon} > 0$ such that for every $A \in C(X, \mathrm{SL}(2, \mathbb{R}))$ with $\|A\|_\infty < C$ and every $\varepsilon \in (0, \bar{\varepsilon}]$, the following holds: Every continuous family of cocycles A_t , $t \in T$ taking values in $\mathcal{B}_\varepsilon(A)$ is contractible in $\mathcal{B}_\varepsilon(A)$.*

Proof. Easy and left to the reader. \square

Now we are able to prove the second main result of this section:

Proof of Proposition 9.2. Let A_t , $t \in \partial[0,1]^p$ be a family of cocycles that is contractible in \mathcal{B} , such that no A_t is uniformly hyperbolic.

Let $V_n = [2^{-n}, 1 - 2^{-n}]^p$. Take n_0 large (to be determined later). Most of the proof will be devoted to the description of a procedure to extend the family of cocycles A_t to the set $[0,1]^p \setminus \mathrm{int} V_{n_0}$, which we call the *shell*.

For each $n \geq n_0$, let $W_n = V_{n+1} \setminus \mathrm{int} V_n$. Cover the set W_n by closed cubes $W_{n,i}$ of edge length 2^{-n-1} whose interiors are disjoint. Let \mathcal{P} be the family of all those cubes $W_{n,i}$. Consider the sets $K_0 \subset K_1 \subset \dots \subset K_p$, where K_q is the union of the q -dimensional faces of cubes in \mathcal{P} . We are going to define A_t successively for t in K_0, K_1, \dots , up to K_p , and so obtain the extension to the shell.

During this proof, we take in \mathbb{R}^p the box distance given by $d((t_i), (t'_i)) = \max_i |t_i - t'_i|$. Fix a sequence $\varepsilon_1 > \varepsilon_2 > \dots$ converging to 0 such that if $\tau, \tau' \in \partial[0,1]^p$ are 2^{-n} -close, then $\|A_\tau - A_{\tau'}\|_\infty < \varepsilon_n$.

Take $t \in K_0$ and let us define A_t . First choose some $\tau = \tau(t) \in \partial[0,1]^p$ such that the distance from t to τ is as small as possible; then that distance equals 2^{-n} for some $n \geq 2$. Since A_τ is not uniformly hyperbolic, by Theorem 5 there is a cocycle B with $\|B - A_\tau\| < \varepsilon_n$ that admits a continuous section. Set $A_t = B$.

The definition of A_t proceeds by an inductive procedure. Assume that A_t is already defined for $t \in K_q$. Fix F to be a $(q+1)$ -dimensional face in K_{q+1} , and let $B = B(F)$ be its relative boundary, that is, the union of the q -dimensional faces in K_q that intersect F . Choose any point t_0 in B , and let $r = r(B)$ be the diameter of $\{A_t : t \in B\}$ in $C(X, \mathrm{SL}(2, \mathbb{R}))$. If r is sufficiently small, then, by Lemma 9.3 combined with Proposition 9.1, we can extend the continuous family of cocycles A_t , $t \in B$ (which takes values in the ball $\mathcal{B}_{2r}(A_0)$) and its continuous family of sections to the $(q+1)$ -dimensional cube F , so that the extended family of cocycles takes values in $\mathcal{B}_{4r}(A_0)$. The smallness of r will be proven by induction, of course. In fact, we will prove that there exist small numbers $r_n^{(q)} > 0$ with $\lim_{n \rightarrow \infty} r_n^{(q)} = 0$ such that if F as above has diameter 2^{-n} , then $r(B) \leq r_n^{(q)}$.

First consider $q = 0$, so F is an edge and $B = \{t_0, t_1\}$. If 2^{-n} is the edge length, then t_0 and t_1 are both within distance 2^{-n+1} of the boundary of the unit cube, hence $d(t_i, \tau(t_i)) \leq 2^{-n+1}$. It follows that $d(\tau(t_0), \tau(t_1)) < 2^{-n+3}$ and so

$$\begin{aligned} r(B) &= \|A_{t_0} - A_{t_1}\|_\infty \leq \|A_{t_0} - A_{\tau(t_0)}\|_\infty + \|A_{\tau(t_0)} - A_{\tau(t_1)}\|_\infty + \|A_{\tau(t_1)} - A_{t_1}\|_\infty \\ &\leq \varepsilon_{n-1} + \varepsilon_{n-3} + \varepsilon_{n-1} < 3\varepsilon_{n-3}. \end{aligned}$$

Thus we set $r_n^{(0)} = 3\varepsilon_{n-3}$. These numbers will all be small for all $n \geq n_0$, provided n_0 is chosen large enough. This proves the $q = 0$ case.

Now let us explain how to find $r_n^{(q)}$ assuming $r_n^{(q-1)}$ is known. Take a $(q+1)$ -dimensional cell F ; its relative boundary $B = B(F)$ is composed of 2^{q+1} cells of dimension q . In each of these cells, the variation of A_t is at most $4r_n^{(q-1)}$. Hence a crude estimate gives $r(B) \leq 2^{q+3}r_n^{(q-1)} =: r_n^{(q)}$. Take n_0 large enough so that these numbers will be small for all $n \geq n_0$.

At this point we have defined A_t for t in the shell. To check continuity in the boundary of the unit cube, take a sequence t_j in $(0, 1)^p \setminus V_{n_0}$ converging to some $\tau \in \partial[0, 1]^p$. For each j , consider a cube of the family \mathcal{P} that contains t_j , and let t'_j be one of its vertices. Then $\|A_{t_j} - A_{t'_j}\|_\infty \rightarrow 0$ as $j \rightarrow \infty$ due to the variation estimates. On the other hand, $\tau(t'_j) \rightarrow \tau$ and $\|A_{\tau(t'_j)} - A_{t'_j}\|_\infty \rightarrow 0$, therefore $A_{t_j} \rightarrow A_\tau$.

Now, taking a smaller shell if necessary, we can assume that all the values of A_t belong to the open set \mathcal{B} . Finally, apply Proposition 9.1 once more to extend continuously the family of cocycles from the shell to the whole unit cube, still taking values in \mathcal{B} , and in such a way that in $(0, 1)^p$, there is a continuous family of sections. This concludes the proof of Proposition 9.2. \square

10. PROOF OF THEOREMS 6 AND 7

Let $A \in C_0(X, \mathrm{SL}(2, \mathbb{R})) \setminus \mathcal{UH}$ be an unlocked cocycle. Fix a continuous determination of the fibered rotation number $\rho : \mathcal{B}_{\varepsilon_0}(A) \rightarrow \mathbb{R}$ on the ball of radius ε_0 and center A . As proved in Appendix C (see Proposition C.1), we have

$$\rho(R_{-\theta}A) < \rho(A) < \rho(R_\theta A) \quad \text{for every sufficiently small } \theta > 0.$$

The first step is to show the following:

Lemma 10.1. *There is a continuous family of cocycles A_t , $t \in [-1, 1]$ such that $A_0 = A$,*

$$\rho(A_{-t}) < \rho(A) < \rho(A_t) \quad \text{for every } t \in (0, 1],$$

and no A_t is uniformly hyperbolic.

Proof. Choose $\theta_n < 0$ an increasing sequence close to 0 and converging to 0, such that $2\rho(R_{\theta_n}A) \notin G$ for each n . The sequence $\rho(R_{\theta_n}A)$ is also strictly increasing (cf. Proposition C.1). So, by continuity of ρ , we can find a sequence δ_n decreasing to 0 such that

$$\rho(R_{\theta_{n-1}}A') < \rho(R_{\theta_n}A'') < \rho(R_{\theta_{n+1}}A''') < \rho(A''') \text{ for all } A', A'', A''', A''' \in \mathcal{B}_{\delta_n}(A).$$

Define an open subset of $C(X, \mathrm{SL}(2, \mathbb{R}))$ by:

$$\mathcal{T}_n = \bigcup_{\theta \in [\theta_n, \theta_{n+1}]} R_\theta \mathcal{B}_{\delta_n}(A).$$

Claim. For each $B \in \mathcal{T}_n$, there exists a unique $\beta \in [-\theta_{n+1} + \theta_{n-1}, -\theta_n]$ such that $\rho(R_\beta B) = \rho(R_{\theta_{n+1}}A)$. Moreover, the function $B \mapsto \beta$ is continuous.

Proof. Let $B \in \mathcal{T}_n$, so $B = R_\alpha A'$ for some $\alpha \in [\theta_n, \theta_{n+1}]$ and $A' \in \mathcal{B}_{\delta_n}(A)$. Then

$$\rho(R_{-\theta_{n+1} + \theta_{n-1}}B) \leq \rho(R_{\theta_{n-1}}A') < \rho(R_{\theta_{n+1}}A) < \rho(A') \leq \rho(R_{-\theta_n}B).$$

By the Intermediate Value Theorem, there is $\beta \in [-\theta_{n+1} + \theta_{n-1}, -\theta_n]$ such that $\rho(R_\beta B) = \rho(R_{\theta_{n+1}}A)$. Uniqueness is a consequence of Proposition C.1 and the fact that $2\rho(R_{\theta_{n+1}}A) \notin G$. Continuity follows. \square

For each n , consider the family of cocycles $R_\theta A$, $\theta \in \{\theta_{n+1}, \theta_n\}$. It is obviously contractible in \mathcal{T}_n . Thus we can apply Proposition 9.2 to find a continuous one-parameter family of cocycles with values in $\mathcal{T}_n \cap \mathcal{UH}^c$, joining $R_{\theta_n} A$ to $R_{\theta_{n+1}} A$. By concatenation, we find a continuous family of cocycles A_t , $t \in [\theta_1, 0]$ such that

$$\begin{aligned} A_{\theta_n} &= R_{\theta_n} A, \\ t \in [\theta_n, \theta_{n+1}] &\Rightarrow A_t \in \mathcal{T}_n \cap \mathcal{UH}^c. \end{aligned}$$

Let

$$t_n = \inf \{t \in [\theta_n, \theta_{n+1}] : \rho(A_{n,t}) = \rho(R_{\theta_{n+1}} A)\}.$$

Now let us define a ‘‘correction’’ function $\beta : [\theta_1, 0] \rightarrow \mathbb{R}$. Define $\beta(t) = 0$ for $t \in [\theta_n, t_n]$. For each $t \in [t_n, \theta_{n+1}]$, let $\beta(t)$ be given by the claim above such that $\rho(R_{\beta(t)} A_{n,t}) = \rho(R_{\theta_{n+1}} A)$. Finally, put $\beta(0) = 0$. Then β is well-defined and continuous.

Now consider the continuous family of cocycles $\tilde{A}_t = R_{\beta(t)} A_t$. We have $\tilde{A}_t \in \mathcal{UH}^c$ for each t (indeed, for $t \in [t_n, \theta_{n+1}]$, this follows from $2\rho(\tilde{A}_t) \notin G$). Also, in each interval $[\theta_n, \theta_{n+1}]$ we have $\rho(\tilde{A}_t) \leq \rho(R_{\theta_{n+1}} A)$, so $\rho(\tilde{A}_t) < \rho(A)$ for every $t < 0$.

Applying an entirely analogous procedure, we find a family \tilde{A}_t for small non-negative t , such that $A_0 = A$ and $\rho(\tilde{A}_t) > \rho(A)$ for $t > 0$. By reparametrizing, we find the desired family. \square

Next we prove the following strengthening of Lemma 10.1:

Lemma 10.2. *There is a continuous family of cocycles A_t , $t \in [-1, 1]$ such that $A_0 = A$,*

$$\rho(A_{-t}) < \rho(A) < \rho(A_t) \quad \text{for every } t \in (0, 1],$$

and the family A_t , $t \in [-1, 1] \setminus \{0\}$ admits a continuous family of sections.

Proof. Let A_t , $t \in [-1, 1]$ be the continuous family of cocycles given by the previous lemma. We can assume it takes values in a small neighborhood of A . Consider the continuous family B_t , $t \in \partial[0, 1]^2$ defined by $B_{(t_1, t_2)} = A_{t_1}$. Since these cocycles are not uniformly hyperbolic, we can apply Proposition 9.2 to find an extended family B_t , $t \in [0, 1]^2$ such that the restricted family B_t , $t \in (0, 1)^2$ is continuously reducible to rotations. Since $\rho(B_{(t_1, 0)}) > \rho(A)$ for each $t_1 \in (0, 1]$, we can find a continuous function $h : (0, 1] \rightarrow (0, 1/2]$ such that $\rho(B_{(t_1/2, h(t_1))}) > \rho(A)$ for each t_1 . That is, $\tilde{A}_t = B_{(t/2, h(t))}$, $t \in (0, 1]$ defines a continuous family of cocycles that admits a continuous family of sections, satisfies $\rho(\tilde{A}_t) > \rho(A)$ for each t , and $\lim_{t \rightarrow 0} \tilde{A}_t = A$. In the same way we find the desired cocycles for negative parameter. \square

In the next step we find a family depending on two parameters:

Lemma 10.3. *There exists a family A_t , $t \in [0, 1]^2$ such that $A_{(0,0)} = A$,*

$$\rho(A_{(0,s)}) < \rho(A) < \rho(A_{(s,0)}) \quad \text{for } s \in (0, 1],$$

and the restricted family A_t , $t \in [0, 1]^2 \setminus \{(0, 0)\}$ admits a continuous family of sections.

Proof. Let A'_t , $t \in [-1, 1]$ be given by Lemma 10.2, so $\rho(A'_{-t}) < \rho(A) < \rho(A'_t)$ for $t > 0$, and A'_t , $t \in [-1, 1] \setminus \{0\}$ admits a continuous family of sections z'_t .

By Proposition 9.1, we can find a family A''_t , $t \in [-1, 1]$, admitting a continuous family of sections z''_t , such that $A''_{\pm 1} = A'_{\pm 1}$ and $z''_{\pm 1} = z'_{\pm 1}$. Next define A_t for $t \in \partial[0, 1]^2$ by

$$A_{(t_1, t_2)} = \begin{cases} A'_{t_1 - t_2} & \text{if } t_1 = 0 \text{ or } t_2 = 0, \\ A''_{t_1 - t_2} & \text{if } t_1 = 1 \text{ or } t_2 = 1. \end{cases}$$

Define analogously the family of sections z_t , $t \in \partial[0, 1]^2 \setminus \{(0, 0)\}$.

Consider the set $\Gamma = \bigcup_{n=1}^{\infty} \partial[0, 1/n]^2$. Using Proposition 9.1 (with $p = 1$) infinitely many times, find an extended continuous family of cocycles A_t , $t \in \Gamma$ such that the cocycles A_t , $t \in \Gamma \setminus \{(0, 0)\}$ admit a continuous family of sections that extends the previously defined z_t .

Now consider the regions $[0, 1/n]^2 \setminus (0, 1/(n+1))^2$; each is homeomorphic to a square, and A_t , B_t are defined in the boundary. Applying Proposition 9.1 (with $p = 2$) we extend A_t , z_t to each of those holes. This can be done so that $\lim_{t \rightarrow (0, 0)} A_t = A$. This gives the desired family of cocycles. \square

Now we need a purely topological result:

Lemma 10.4. *Let $\theta : [0, 1]^2 \rightarrow \mathbb{R}$ be a continuous function. Assume that $\theta(0, 0) = 0$ and $\theta(0, s) < 0 < \theta(s, 0)$ for every $0 < s \leq 1$. Let $K = \theta^{-1}(0)$ and let $U \subset (0, 1]^2$ be a neighborhood of $K \setminus \{(0, 0)\}$. Then there exists a path $\gamma : (0, 1] \rightarrow U$ such that $\lim_{s \rightarrow 0} \gamma(s) = (0, 0)$.*

Proof. Let $\varepsilon : (0, 1] \rightarrow (0, 1]$ be a continuous increasing function such that for every $t \in K \setminus \{(0, 0)\}$, $B_{\varepsilon(|t|)}(t) \cap [0, 1]^2 \subset U$.

For every $t \in K \setminus \{(0, 0)\}$, choose a closed dyadic square \tilde{D}_t containing t of diameter at least $\varepsilon(|t|)/10$ and at most $\varepsilon(|t|)/5$, and let D_t be the union of all closed dyadic squares of the same diameter that intersect \tilde{D}_t . Let $W = \bigcup_{t \in K \setminus \{0\}} D_t$.

Let $I = \{t = (t_1, t_2) \in [0, 1]^2; t_1 + t_2 = 1/2\}$. Then I intersects W on finitely many of its connected components, which we denote by S_i . We claim that one of them accumulates at $(0, 0)$. Indeed, if this is not the case, let T be the union of the boundaries of the S_i intersected with $\Delta = \{t = (t_1, t_2) \in [0, 1]^2 : t_1 + t_2 \leq 1/2\}$. Then there exists a path joining $(0, 1/2)$ to $(1/2, 0)$ that is contained everywhere in $(I \setminus \bigcup S_i) \cup T$. But the sign of θ is constant along this path, a contradiction.

Let S be such a connected component that accumulates at $(0, 0)$. Write S as a countable union of dyadic squares, so that any of them does not contain another. Order such squares with the property that $D_i \cap D_{i+1} \neq \emptyset$. Refine this sequence so that no squares are repeated. Join the centers of the squares by straight segments. This gives the desired path. \square

Now we can finish the proof of Theorems 6 and 7:

Proof. Consider the continuous family of cocycles \tilde{A}_t , $t \in [0, 1]^2$ given by Lemma 10.3; for $t \neq (0, 0)$, one can write

$$\tilde{A}_t(x) = \tilde{B}_t(fx) R_{2\pi\theta(t,x)} \tilde{B}_t(x)^{-1},$$

where $\tilde{B}_t \in C(X, \mathrm{SL}(2, \mathbb{R}))$, $t \in [0, 1]^2 \setminus \{(0, 0)\}$ is a continuous family of conjugacies and $\theta : [0, 1]^2 \times X \rightarrow \mathbb{R}$ is continuous, with $\int \theta(0, s, x) d\mu(x) < \int \theta(0, 0, x) d\mu(x) < \theta(s, 0, x)$ for every $0 < s \leq 1$. Here, we write $\int \theta(0, 0, x) d\mu(x)$ for $\lim_{t \rightarrow (0, 0)} \int \theta(t, x) d\mu(x)$. By definition of the fibered rotation number, $\int \theta(0, 0, x) d\mu(x) = \rho(A) \bmod G'$ (where G' is a subgroup of G , see §1.4).

Let $g = \alpha - \int \theta(0, 0, x) d\mu(x)$, in the case of Theorem 6, and let $g = 2\alpha - 2 \int \theta(0, 0, x) d\mu(x)$, in the case of Theorem 7. In either case, $g \in G$, so we may write $g = \int \phi d\mu$ for some continuous function $\phi : X \rightarrow \mathbb{R}$ such that $\phi = \psi \circ f - \psi$ for some continuous $\psi : X \rightarrow \mathbb{R}/\mathbb{Z}$. In the case of Theorem 6, we may thus replace $\tilde{B}_t(x)$ by $\tilde{B}_t(x) R_{2\pi\psi(x)}$ and $\theta(t, x)$ by $\theta(t, x) + \phi(t, x)$, that is, we may actually assume that $g = 0$. In the case of Theorem 7, we may proceed similarly, replacing $\tilde{B}_t(x)$ by the $\mathrm{PSL}(2, \mathbb{R})$ -valued $\tilde{B}_t(x) R_{\psi(x)}$, and $\theta(t, x)$ by $\theta(t, x) + \phi(t, x)/2$: we may thus also assume that $g = 0$ in this case. From now on, we assume that $\alpha = \int \theta(0, 0, x) d\mu(x)$.

Apply Lemma 7.1 with $T = [0, 1]^2 \setminus \{(0, 0)\}$, $\varphi = \theta|(T \times X)$, $T^* = \emptyset$ and ε so small that $\varepsilon(0, s) < \int \alpha - \theta(0, s, x) d\mu(x)$, $\varepsilon(s, 0) < \int \theta(s, 0, x) - \alpha d\mu(x)$ and

$\|\tilde{B}_t\|_\infty^2 \varepsilon(t) \rightarrow 0$ as $t \rightarrow (0, 0)$. Let $\tilde{\varphi}$ and $\tilde{\psi}$ be as in the lemma. We now define $\hat{A}_t(x) = \tilde{B}_t(f(x))R_{2\pi\tilde{\varphi}(t,x)}\tilde{B}_t(x)^{-1}$ for $t \in [0, 1]^2 \setminus \{(0, 0)\}$ and $\hat{A}_0 = A$: this is a continuous family at $t = (0, 0)$ since $\|\hat{A}_t - \tilde{A}_t\|_\infty < 2\pi\varepsilon(t)\|\tilde{B}_t\|_\infty^2$ for $t \in [0, 1]^2 \setminus \{(0, 0)\}$. Since $\tilde{\varphi}(t, x)$ is cohomologous to $c(t) = \int \theta(t, x) d\mu(x) = \rho(\tilde{A}_t)$, we can write $\hat{A}_t(x) = \tilde{B}_t(f(x))R_{2\pi c(t)}\tilde{B}_t(x)^{-1}$ (where the ψ 's are absorbed by \hat{B}).

Now choose some function $\varepsilon_1(t) > 0$ that goes to zero as $t \rightarrow (0, 0)$. By Lemma 10.4, there exists a continuous path $\gamma : (0, 1] \rightarrow (0, 1]^2$ whose range is contained in the set of t 's for which $\|\hat{A}_t(x) - \tilde{B}_t(f(x))R_{2\pi\alpha}\tilde{B}_t(x)^{-1}\| < \varepsilon_1(t)$. (The latter set is open and contains $K = c^{-1}(\alpha)$.) Thus, we may set $A_s(x) = \hat{B}_{\gamma(s)}(f(x))R_{2\pi\alpha}\hat{B}_{\gamma(s)}(x)^{-1}$. By construction, this family has the desired properties. \square

Remark 10.5. It follows from the proof that if $\alpha = \rho_{f,A} \pmod{G'}$, then the maps B_t can be taken in $C_0(X, \mathrm{SL}(2, \mathbb{R}))$ (i.e., homotopic to constant). Conversely, if the conclusion of Theorem 6 holds with $B_t \in C_0(X, \mathrm{SL}(2, \mathbb{R}))$, then it is straightforward to show that $\rho_{f,A} = \alpha \pmod{G'}$.

Remark 10.6. A cocycle can be accessed by a continuous path of unlocked cocycles with the same rotation number if and only if it is unlocked itself. Indeed, the “only if” part is an immediate consequence of the definitions. On the other hand, notice that if two cocycles in $C_0(X, \mathrm{SL}(2, \mathbb{R}))$ are conjugate via a map also in $C_0(X, \mathrm{SL}(2, \mathbb{R}))$ then⁹ they have the same type (locked / unlocked / semi-locked). Thus, given an unlocked cocycle A , applying Theorem 6 (and Remark 10.5) with $\alpha = \rho_{f,A}$ we obtain a continuous path of unlocked cocycles with the same rotation number.

APPENDIX A. CONNECTEDNESS CONSIDERATIONS

Recall that $C_0(X, \mathrm{SL}(2, \mathbb{R}))$ indicates the subset of $C(X, \mathrm{SL}(2, \mathbb{R}))$ formed by maps that are homotopic to a constant. The rotation number can be viewed as a map $C_0(X, \mathrm{SL}(2, \mathbb{R})) \rightarrow \mathbb{R}/G'$.

Proposition A.1. *If two cocycles have the same rotation number ($\pmod{G'}$), then is it possible to connect them by a continuous path with constant rotation number ($\pmod{G'}$).*

Proof. Let $A', A'' \in C_0(X, \mathrm{SL}(2, \mathbb{R}))$ be two cocycles such that $\rho_{f,A'} = \rho_{f,A''} \pmod{G'}$. Fix once and for all a lift $F' : X \times \mathbb{R} \rightarrow X \times \mathbb{R}$ such that $e^{2\pi i F'_2(x,t)}$ is a positive real multiple of $A'(x) \cdot e^{2\pi i t}$. This determines the fibered rotation number of (f, A') as some $a \in \mathbb{R}$. Moreover, for any contractible space Y and any continuous map $y \in Y \mapsto A_y \in C_0(x, \mathrm{SL}(2, \mathbb{R}))$ such that $A_{y_0} = A'$ for some $y_0 \in Y$, there exists a unique continuous function $\rho : Y \rightarrow \mathbb{R}$, such that $\rho(y) = \rho_{f,A_y} \pmod{G'}$ and $\rho(y_0) = a$. We also call ρ a continuous determination of the rotation number. (In fact, the previous use of this terminology is a particular case.)

Take any continuous path $t \in [0, 1] \mapsto A_t \in C_0(X, \mathrm{SL}(2, \mathbb{R}))$ with $A_0 = A'$ and $A_1 = A''$. Consider the continuous determination of the rotation number $\rho : [0, 1] \rightarrow \mathbb{R}$ with $\rho(0) = a$. Then $\rho(1) - \rho(0) \in G'$, thus there exists $\phi \in C(X, \mathbb{Z})$ such that $\rho(1) - \rho(0) = \int \phi d\mu$. Extend the path of cocycles to the interval $[0, 2]$ by defining $A_t(x) = R_{-2\pi(t-1)\phi(x)}A''(x)$ for $t \in [1, 2]$. Considering the corresponding extension $\rho : [0, 2] \rightarrow \mathbb{R}$, we have $\rho(2) = \rho(0) = a$.

Now consider the two-parameter family $(t, \theta) \in [0, 2] \times \mathbb{R} \mapsto R_\theta A_t \in C_0(X, \mathrm{SL}(2, \mathbb{R}))$ and the associated continuous determination of the rotation number $\rho : [0, 2] \times \mathbb{R} \rightarrow \mathbb{R}$ such that $\rho(0, 0) = a$, and thus $\rho(0, 2) = a$. Notice that $\rho(t, \theta)$ is non-decreasing as a function of θ , and $\rho(t, \theta + 2\pi) = \rho(t, \theta) + 1$.

⁹In fact, it is not necessary that the conjugacy is homotopic to constant; see Proposition C.4.

We claim that the function $f(t) = \inf\{\theta \in \mathbb{R} : \rho(t, \theta) = a\}$ is continuous. Indeed, fix any sequence t_n in $[0, 2]$ converging to some t and such that $\beta = \lim f(t_n)$ exists. Let $\alpha = f(t)$. Then $\rho(t, \beta) = \lim \rho(t_n, \beta) \leq a = \rho(t, \alpha)$. It follows that $\beta \leq \alpha$. Moreover, if this inequality is strict, then $\rho(t, \theta)$ is constant equal to a for θ in the interval (β, α) and thus by the first part of Proposition C.1, $R_{(\alpha+\beta)/2}A_t$ is uniformly hyperbolic. In particular $\rho(s, (\alpha+\beta)/2) = a$ for every s sufficiently close to t . This implies that $\rho(t, \beta) \leq \lim \rho(t_n, \beta) \leq \lim \rho(t_n, (\alpha+\beta)/2, t_n) = a$ and thus, $\beta \geq f(t)$ – contradiction.

It now follows that the rotation number is constant mod G' along the continuous path of cocycles $t \in [0, 2] \mapsto R_{\min(0, f(t))}A_t$, concluding the proof. \square

Proposition A.2. *If $A_0, A_1 \in C_0(X, \mathrm{SL}(2, \mathbb{R})) \setminus \mathcal{UH}$, then we can join these cocycles by a continuous path $t \in [0, 1] \mapsto A_t \in C(X, \mathrm{SL}(2, \mathbb{R})) \setminus \mathcal{UH}$. Moreover, if A_0 and A_1 have the same rotation number (mod G'), then we can take the path with constant rotation number (mod G').*

Proof. The first assertion is a trivial consequence of Proposition 9.2 with $p = 1$.

To prove the second assertion, it suffices to show that for any cocycle $A \in C_0(X, \mathrm{SL}(2, \mathbb{R})) \setminus \mathcal{UH}$ and any $\alpha = \rho_{f, A}$ mod G , there is a path joining A and $R_{2\pi\alpha}$ along which the rotation number is constant.

Given such A and α , apply Theorem 6 and find continuous paths A_t , $t \in [0, 1]$ and B_t , $t \in (0, 1]$ in $C(X, \mathrm{SL}(2, \mathbb{R}))$ such that $A_0 = A$ and $A_t(x) = B_t(f(x))R_{2\pi\alpha}B_t(x)^{-1}$.

Since $\mathrm{SO}(2, \mathbb{R})$ is a deformation retract of $\mathrm{SL}(2, \mathbb{R})$, there exists a homotopy $h_t : \mathrm{SL}(2, \mathbb{R}) \rightarrow \mathrm{SL}(2, \mathbb{R})$, $t \in [0, 1]$ between the identity h_0 and a $\mathrm{SO}(2, \mathbb{R})$ -valued map h_1 .

The rotation number is constant along the path $A_{t+1}(x) = h_t(B_1(fx))R_{2\pi\alpha}[h_t(B_1(x))]^{-1}$, $t \in [0, 1]$ joining A_1 and the $\mathrm{SO}(2, \mathbb{R})$ -valued cocycle A_2 . The latter is homotopic to constant as a map $X \rightarrow \mathrm{SL}(2, \mathbb{R})$ and thus (using the deformation retract) also as a map $X \rightarrow \mathrm{SO}(2, \mathbb{R})$. It follows that we can write $A_2(x) = R_{\theta_0(x)}$ where θ_0 is real-valued. Let θ_t , $t \in [0, 1]$ be a homotopy in $C(X, \mathbb{R})$ between θ_0 and constant $2\pi\alpha$ such that $\int \theta_t d\mu$ does not depend on t . Define $A_{2+t}(x) = R_{\theta_t(x)}$ for $t \in [0, 1]$. Then A_t , $t \in [0, 3]$ is a path joining A and $R_{2\pi\alpha}$ along which the rotation number is constant. \square

On the other hand, the set of uniformly hyperbolic cocycles with a given rotation number is not necessarily connected:

Proposition A.3. *There exists a strictly ergodic homeomorphism f of a finite-dimensional compact space X and there exist two cocycles $A_0, A_1 \in C_0(X, \mathrm{SL}(2, \mathbb{R}))$ with the same rotation number that lie in distinct path-connected components of \mathcal{UH} .*

Remark A.4. One could raise the question of whether it is possible to add an extra parameter in Theorem 3 and obtain the following statement: “If A_s , $s \in [0, 1]$ is a path in $C_0(X, \mathrm{SL}(2, \mathbb{R}))$ with $2\rho_{f, A_s}$ constant and in G , then there is a continuous family of cocycles $A_{s,t}$, $(s, t) \in [0, 1]^2$ such that $A_{s,0} = A_s$ and $A_{s,t} \in \mathcal{UH}$ for $t > 0$.” However, it follows from Propositions A.1 and A.3 that such a statement is *false* in general. (On the other hand, it is possible to add an extra parameter in Theorem 4: we have done that already in Proposition 9.2.)

To prove Proposition A.3 we will need the following well-known lemma. Lacking a suitable reference, we give a proof of it.

Lemma A.5. *Let $\alpha \in \mathbb{R} \setminus \mathbb{Q}$ and $\beta \in \mathbb{R}$. For a generic continuous $\phi : \mathbb{R}/\mathbb{Z} \rightarrow \mathbb{R}$ with $\int \phi(x) dx = \beta$, the map¹⁰ $f_{\alpha,\phi} : \mathbb{R}^2/\mathbb{Z}^2 \rightarrow \mathbb{R}^2/\mathbb{Z}^2$ given by $f_{\alpha,\phi}(x, y) = (x + \alpha, y + \phi(x))$ is strictly ergodic.¹¹*

Proof. Let $h_i : \mathbb{R}^2/\mathbb{Z}^2 \rightarrow \mathbb{R}$, $i \geq 1$ be a dense sequence in $C(\mathbb{R}^2/\mathbb{Z}^2, \mathbb{R})$. Let U_n be the set of all continuous ϕ with average β such that for every ergodic invariant measure μ for $f_{\alpha,\phi}$, we have $|\int h_n d\mu - \int h_n d\text{Leb}| < 1/n$.

Clearly, if $\phi \in \bigcap U_n$, then ϕ is uniquely ergodic. So we just have to show that all U_n are dense. Let ϕ_0 be any trigonometric polynomial. We wish to find ϕ close to ϕ_0 that belongs to U_n . Now, clearly $\phi_0(x) = \psi_0(x + \alpha) - \psi_0(x) + \beta$ for some trigonometric polynomial.

For an arbitrary continuous function $g : \mathbb{R}/\mathbb{Z} \rightarrow \mathbb{R}$, let μ_g be the probability measure on $\mathbb{R}^2/\mathbb{Z}^2$ that projects to Lebesgue measure in the first factor and is supported on the graph of $x \mapsto x + g(x)$. Consider the sequence of functions $\delta_k : \mathbb{R}/\mathbb{Z} \rightarrow \mathbb{R}$, such that $\delta_k(x)$ is twice the distance between kx and \mathbb{Z} . It is readily seen that μ_{δ_k} converges in weak-* sense to Leb. It follows that for every continuous function g , $\mu_{g+\delta_k}$ converges to Leb (which equals the push-forward of Leb under $(x, y) \mapsto (x, y + g(x))$).

In particular, $\mu_{\psi_0+\delta_k} \rightarrow \text{Leb}$. Let p/q be a continued fraction approximant of α and take $\psi = \psi_0 + \delta_q$. Let $\phi(x) = \psi(x + \alpha) - \psi(x)$. Then any invariant measure for $f_{\alpha,\phi}$ is a vertical translate of μ_ψ , that is, it is of the form $(T_t)_* \mu_\psi$ where $T_t(x, t) = (x, y + t)$. It follows that all of them are close to Leb, so if we take q sufficiently large, we will have $\phi \in U_n$. Finally, one computes $\|\phi - \phi_0\|_\infty \leq 2|q\alpha - p| < 1/q$, so taking q large also ensures that ϕ is close to ϕ_0 . \square

Proof of Proposition A.3. By the previous lemma, there exists a strictly ergodic homeomorphism $f = f_{\alpha,\phi}$ of the torus $\mathbb{R}^2/\mathbb{Z}^2$ with $\int \phi(x_1) dx_1 = \alpha \in \mathbb{R} \setminus \mathbb{Q}$. Let $\psi : \mathbb{R}/\mathbb{Z} \rightarrow \mathbb{R}$ be continuous and such that $\psi(x_1 + \alpha) - \psi(x_1)$ is uniformly close to $\phi(x_1) - \alpha$. (This can be done easily using the remark after Lemma 3 in [ABD].) Let $u : \mathbb{R}^2/\mathbb{Z}^2 \rightarrow \mathbb{R}/\mathbb{Z}$ be given by:

$$u(x_1, x_2) = x_1 - x_2 + \psi(x_1).$$

Thus $u(f(x_1, x_2)) - u(x_1, x_2) = \psi(x_1 + \alpha) - \psi(x_1) + \alpha - \phi(x_1)$, which is close to 0. Notice that u is not homotopic to constant. Define two elements of $C(\mathbb{R}^2/\mathbb{Z}^2, \text{SL}(2, \mathbb{R}))$ by

$$A_0(x_1, x_2) = \begin{pmatrix} 2 & 0 \\ 0 & 1/2 \end{pmatrix}, \quad A_1(x_1, x_2) = R_{u(f(x_1, x_2))} \begin{pmatrix} 2 & 0 \\ 0 & 1/2 \end{pmatrix} R_{-u(x_1, x_2)}.$$

These cocycles are uniformly hyperbolic and have the same rotation number (because $u \circ f - u$ lifts to a function of zero mean). But no path A_t , $t \in [0, 1]$ joining A_0 and A_1 can be wholly contained in \mathcal{UH} , because the respective expanding directions $e^u(A_0)$, $e^u(A_1) : \mathbb{R}^2/\mathbb{Z}^2 \rightarrow \mathbb{P}^1$ are not homotopic. \square

APPENDIX B. MINIMAL EXAMPLES WITH MANY PERSISTENTLY CLOSED GAPS

In [Be2], Question 4, the question of genericity of Cantor spectrum is asked in a more general context than strictly ergodic dynamical systems. The Gap Labelling Theorem holds, in fact, for any homeomorphism of a compact metric space, and Bellissard asked whether the presence of periodic orbits could be the only obstruction for the genericity of Cantor spectrum (when it is allowed at all by the Gap Labeling Theorem). Though we have answered this question affirmatively in our more restricted setting, the answer to the original question is in fact negative, and

¹⁰Maps of such form are called *Anzai skew-products* after [An].

¹¹ It is well-known that ergodic Anzai skew-products are in fact strictly ergodic; see [Fu], but this will not be used below.

as we will see even a slight departure from unique ergodicity (keeping minimality) may already prevent Cantor spectrum.

We will consider the well-known Furstenberg examples (see [Fu, page 585] or [Pa, § 5.5]) of minimal, but non-uniquely ergodic skew-products of the form $f(x, y) = (x + \alpha, y + \phi(x))$, $\phi \in C(X = \mathbb{R}^2/\mathbb{Z}^2, \mathbb{R})$, $\alpha \in \mathbb{R} \setminus \mathbb{Q}$. For these examples, $G = \mathbb{Z} \oplus \alpha\mathbb{Z}$ is independent of the choice of the invariant measure.

Since \mathcal{L} is dense in $(0, 1)$, the Gap Labelling Theorem does not prevent Cantor spectrum for a Furstenberg example, and in fact there are choices of the sampling function exhibiting Cantor spectrum.¹² However, the next theorem shows that Cantor spectrum (and hence the converse of the Gap Labelling Theorem) is not generic for such dynamics.

Theorem 9. *Let f be a Furstenberg example.*

- a) *For an open and dense set of $v \in C(X, \mathbb{R})$, $\Sigma_{f,v}$ is not a Cantor set,*
- b) *For a generic set of $v \in C(X, \mathbb{R})$, $\text{int } \Sigma_{f,v}$ is dense in $\Sigma_{f,v}$.*¹³
- c) *The interior of $C_0(X, \text{SL}(2, \mathbb{R})) \setminus \mathcal{UH}$ is dense in $C_0(X, \text{SL}(2, \mathbb{R})) \setminus \mathcal{UH}$.*

Proof. We start with the third statement. Notice that if $A \in C_0(X, \text{SL}(2, \mathbb{R})) \cap \overline{\mathcal{UH}}$, the fibered rotation number (viewed as an element of $\mathbb{R}/G' = \mathbb{R}/\mathbb{Z}$) is independent of the invariant measure,¹⁴.

If $A \in C_0(X, \text{SL}(2, \mathbb{R})) \cap C(X, \text{SO}(2, \mathbb{R}))$, it has the form $A(x) = R_{\phi(x)}$ with $\phi : X \rightarrow \mathbb{R}$ continuous, and the fibered rotation number is independent of the invariant measure if and only if $\int_X \phi \, d\mu$ is. Thus for an open and dense set of $A \in C_0(X, \text{SL}(2, \mathbb{R})) \cap C(X, \text{SO}(2, \mathbb{R}))$, $A \notin \overline{\mathcal{UH}}$. By Theorem 5, if A is not uniformly hyperbolic, it can be approximated by a cocycle which is bounded, and hence by one in the complement of \mathcal{UH} . This proves (c).

Using Lemma 1.3, (c) implies that for $v_0 \in C(X, \mathbb{R})$ and $E \in \Sigma_{f,v_0}$, since $A_{E,v_0} \notin \mathcal{UH}$, we can perturb v_0 to v such that $A_{E,v} \notin \overline{\mathcal{UH}}$. Then $E \in \text{int } \Sigma_{f,\tilde{v}}$ for every \tilde{v} in a neighborhood of v . This gives (a). A simple Baire category argument allows us to conclude (b). \square

APPENDIX C. UNIFORM HYPERBOLICITY AND THE ROTATION NUMBER

In this section we discuss in more detail the relationship between the behavior of the rotation number in a neighborhood of a cocycle that is homotopic to a constant and properties of the cocycle. Recall the definitions of a locked/semi-locked/unlocked cocycle, which were given in Section 2.

Proposition C.1. *Let $A \in C(X, \text{SL}(2, \mathbb{R}))$ be homotopic to a constant, and fix a continuous determination of the fibered rotation number $\rho : \mathcal{V} \rightarrow \mathbb{R}$ on a neighborhood of A . Then:*

- a) *The function $r_A : \theta \mapsto \rho(R_\theta A)$, defined on an open interval around 0, is locally constant at $\theta = 0$ if and only if $A \in \mathcal{UH}$. Consequently, the cocycle A is:
 - locked if R_A is locally constant at $\theta = 0$;
 - semi-locked if $\theta = 0$ is in the boundary of a plateau of r_A ;
 - unlocked if the value $\rho(A)$ does not correspond to a plateau of r_A .*
- b) *The cocycle A is uniformly hyperbolic if and only if it is locked.*

¹²Indeed, the converse of the Gap Labelling Theorem holds for generic potentials which depend only on the first coordinate: this case reduces to the strictly ergodic dynamics $x \mapsto x + \alpha$ which is covered by Corollary 2.

¹³It is reasonable to expect that for generic v , $\partial \Sigma_{f,v}$ is a Cantor set (thus $\Sigma_{f,v}$ is not generically a finite union of intervals).

¹⁴It is enough to check this for $A \in C_0(X, \text{SL}(2, \mathbb{R})) \cap \mathcal{UH}$: consider $\rho(A)$ as a continuous function of the convex set of invariant measures, which (since $A \in \mathcal{UH}$) takes values in $\frac{1}{2}\mathbb{Z} + \frac{\alpha}{2}\mathbb{Z}$.

- c) If the cocycle A is semi-locked, it can be accessed by uniformly hyperbolic cocycles of the form $A_\theta = R_\theta A$, where θ runs through a (non-degenerate) interval of the form $[0, \varepsilon]$ or $[-\varepsilon, 0]$.
- d) If A takes values in $\mathrm{SO}(2, \mathbb{R})$, or if $2\rho(A) \notin G$, then A is unlocked.

Before giving the proof of this proposition, let us first establish the following lemma.

Lemma C.2. *For any $\varepsilon > 0$, there exist $c > 0$ and $\lambda > 1$ with the following properties. Let v be a non-zero vector in \mathbb{R}^2 . Assume A_0, A_1, \dots, A_{n-1} are matrices in $\mathrm{SL}(2, \mathbb{R})$ such that*

$$\|A_{n-1} \cdots A_1 A_0(v)\| < c\lambda^n \|v\|.$$

Then for any $w \in \mathbb{R}^2 \setminus \{0\}$, there exists $\theta \in [-\varepsilon, \varepsilon]$ such that the vectors

$$R_\theta A_{n-1} R_\theta A_{n-2} \cdots R_\theta A_0 R_\theta(v) \quad \text{and} \quad w$$

are collinear.

A set $I \subset \mathbb{P}^1$ is called an interval if it is non-empty, connected, and its complement contains more than one point. If $I \subset \mathbb{P}^1$ is an open interval, there is a Riemannian metric $\|\cdot\|_I$ on I , called the *Hilbert metric*, which is uniquely characterized by the following properties:

- any invertible linear map $A : \mathbb{R}^2 \rightarrow \mathbb{R}^2$ induces an isometry between I and $A(I)$ with the corresponding Hilbert metrics;
- if $I \subset \mathbb{P}^1$ is the projectivization of $\{(1, y) \in \mathbb{R}^2 : y \in (0, \infty)\}$, then the Hilbert metric is given by $\frac{dy}{y}$.

If $J \Subset I$ are open intervals, then the metric $\|\cdot\|_I$ restricted to J is smaller than $\|\cdot\|_J$ times a factor $\tanh(\mathrm{diam}_I(J)/4) < 1$ (Garrett Birkhoff's formula).

Proof of Lemma C.2. We can assume that $\varepsilon < \pi/2$, otherwise the conclusion is trivial. Let $|\cdot|$ indicate the Euclidean (angle) metric on \mathbb{P}^1 (so that the length of \mathbb{P}^1 is π). Fix constants $0 < C_1 < C_2$ (depending only on ε) such that for any open interval I whose Euclidean length $|I|$ is between ε and $\pi - \varepsilon$, we have:

$$\begin{aligned} \|\mathbf{u}\|_I &\geq C_1 |\mathbf{u}| \quad \text{for every } \mathbf{u} \text{ tangent to } I; \\ \|\mathbf{u}\|_I &\leq C_2 |\mathbf{u}| \quad \text{if } \mathbf{u} \text{ is a tangent vector at the midpoint of } I. \end{aligned}$$

Also fix a constant $\tau > 1$ such that if $J \Subset I$ are open intervals with the same midpoint and $|J| + \varepsilon = |I| \leq \pi - \varepsilon$, then

$$\|\mathbf{u}\|_J \geq \tau \|\mathbf{u}\|_I \quad \text{for every } \mathbf{u} \text{ tangent to } J.$$

Now take $A_0, \dots, A_{n-1} \in \mathrm{SL}(2, \mathbb{R})$ and a unit vector $v \in \mathbb{R}^2$. Let $\bar{v} \in \mathbb{P}^1$ correspond to v . Assume that

$$(14) \quad \{R_\theta A_{n-1} \cdots A_1 R_\theta A_0 R_\theta(\bar{v}) : \theta \in [-\varepsilon, \varepsilon]\} \neq \mathbb{P}^1;$$

we want to show that $\|A_{n-1} \cdots A_0\|$ is exponentially large.

Let us define some open intervals $I_0, I_1, \dots, I_n \subset \mathbb{P}^1$. Let I_0 be the open interval with midpoint \bar{v} and length $|I_0| = \varepsilon$. Define inductively I_k as the open interval with the same midpoint as $A_{k-1}(I_{k-1})$ and length $|I_k| = |A_{k-1}(I_{k-1})| + \varepsilon$. It follows from (14) and the fact that the matrices A_j preserve orientation that I_k is indeed an interval, with endpoints $R_{-\varepsilon/2} A_{k-1} \cdots R_{-\varepsilon/2} A_0 R_{-\varepsilon/2}(\bar{v})$ and $R_{\varepsilon/2} A_{k-1} \cdots R_{\varepsilon/2} A_0 R_{\varepsilon/2}(\bar{v})$, and length $|I_k| \leq \pi - \varepsilon$. Therefore for any vector \mathbf{u} tangent to I_{k-1} , we have $\|D A_{k-1}(\mathbf{u})\|_{I_k} \leq \tau^{-1} \|\mathbf{u}\|_{I_{k-1}}$.

Let \mathbf{w} be a non-zero tangent vector at \bar{v} , and write $B = A_{n-1} \cdots A_0$. Then

$$|DB(\mathbf{w})| \leq C_1^{-1} \|DB(\mathbf{w})\|_{I_n} \leq C_1^{-1} \tau^{-n} \|\mathbf{w}\|_{I_0} \leq C_1^{-1} C_2 \tau^{-n} |\mathbf{w}|.$$

Since $\|B(v)\|^{-2} = |DB(\mathbf{w})|/|\mathbf{w}|$, ¹⁵ we get that $\|B(v)\|$ is exponentially large, as desired. \square

Proof of Proposition C.1. We begin by proving the first statement in part (a). Obviously, ρ is locally constant around A if it is uniformly hyperbolic. Let us prove the other direction. As usual, fix a lift $F_2 : X \times \mathbb{R} \rightarrow \mathbb{R}$ of A . The nearby lift of $R_\theta A$ is $F_2 + \theta$, hence the rotation number is given by

$$\rho(R_\theta A) = \lim_{n \rightarrow \infty} \frac{1}{n} \int [(F_2 + \theta)^n(x, t) - t] d\mu(x) \quad (t \text{ arbitrary}),$$

where $(F_2 + \theta)^n(x, t)$ is the second coordinate of the n -th iterate of $(x, t) \mapsto (f(x), F_2(x, t) + \theta)$.

Assume that (f, A) is not uniformly hyperbolic. Fix $\varepsilon > 0$. We will show that $\rho(R_\varepsilon A) > \rho(R_{-\varepsilon} A)$. Let $c > 0$ and $\lambda > 1$ be given by Lemma C.2, depending on ε . There exist $n_0 \in \mathbb{N}$ and $x \in X$ such that $\|A^{n_0}(x)\| < c\lambda^{n_0}$. Let G indicate the (positive measure) set of such points x . Then Lemma C.2 implies that

$$(F_2 + \varepsilon)^{n_0}(x, t) - (F_2 - \varepsilon)^{n_0}(x, t) > \frac{1}{2} \quad \text{for all } x \in G, t \in \mathbb{R}.$$

When $x \notin G$, the left-hand side is positive. Since orbits under f^{n_0} visit G with an average frequency $\mu(G)$, we conclude that

$$\rho(R_\varepsilon A) - \rho(R_{-\varepsilon} A) = \lim_{k \rightarrow \infty} \frac{1}{n_0 k} \int [(F_2 + \varepsilon)^{n_0 k}(x, t) - (F_2 - \varepsilon)^{n_0 k}(x, t)] d\mu(x) \geq \frac{\mu(G)}{2n_0},$$

as we wanted to show. This proves the first part of (a).

Let us now prove the remaining parts of (a). Suppose A is such that $r_A : \theta \mapsto \rho(R_\theta A)$ is constant on $[0, \varepsilon]$. By what we just proved, it follows that $R_{\varepsilon/2} A \in \mathcal{UH}$. Take a neighborhood V of A such that $R_{\varepsilon/2} B \in \mathcal{UH}$ and $\rho(R_{\varepsilon/2} B) = \rho(R_{\varepsilon/2} A) = \rho(A)$ for every $B \in V$. Since r_A is non-decreasing, we get $\rho(B) \leq \rho(R_{\varepsilon/2} B) = \rho(A)$ and hence ρ attains a local maximum at A . Analogously, if A is such that $\theta \mapsto \rho(R_\theta A)$ is constant on $[-\varepsilon, 0]$, then ρ attains a local minimum at A . This establishes the remaining three implications in part (a).

The parts (b), (c), and (d) follow readily from (a). (To see the second part of (d), note that if A is not unlocked, then $\theta \mapsto \rho(R_\theta A)$ is constant equal to ρ_0 on a interval I of the form $[0, \varepsilon]$ or $[-\varepsilon, 0]$. Thus, $R_\theta A$ is uniformly hyperbolic for all $\theta \in \text{int } I$ and hence $2\rho(A) \in G$.) \square

Remark C.3. As is apparent from the proof, Proposition C.1 holds in the more general setting where f is a homeomorphism of a compact space with a fully supported invariant probability measure μ and the rotation numbers are taken with respect to this measure (see [H], §5.8.2).

Proposition C.4. *Conjugate cocycles have the same type (locked / semi-locked / unlocked).*

Proof. Let $B \in C(X, \text{SL}(2, \mathbb{R}))$ be such that $x \mapsto B(f(x))B(x)^{-1}$ is homotopic to a constant and let $\Phi : C_0(X, \text{SL}(2, \mathbb{R})) \rightarrow C_0(X, \text{SL}(2, \mathbb{R}))$ be given by $\Phi(A)(x) = B(f(x))A(x)B(x)^{-1}$.

Let \mathcal{U} be an open connected subset of $C_0(X, \text{SL}(2, \mathbb{R}))$ such that there exists a continuous choice of lifts $A \in \mathcal{U} \mapsto F_{A,2} \in C(X \times \mathbb{R}, \mathbb{R})$. These lifts induce determinations $\rho : \mathcal{U} \rightarrow \mathbb{R}$ and $\rho' : \Phi(\mathcal{U}) \rightarrow \mathbb{R}$ of the fibered rotation number. We are going to show that there exists $c \in \mathbb{R}$ such that $\rho'(\Phi(A)) = \rho(\Phi(A)) + c$ for every $A \in \mathcal{U}$. The proposition follows immediately.

¹⁵Indeed, the slice of the unit disk corresponding to a small angle $|\mathbf{w}|$ around v is a region of \mathbb{R}^2 with area $\|v\|^2|\mathbf{w}|$ that is mapped by B to a region of area $\|B(v)\|^2|DB(\mathbf{w})| + \mathcal{O}(|\mathbf{w}|^2)$; but B is area-preserving.

Let also $\Theta : X \times \mathbb{R} \rightarrow \mathbb{R}$ be a measurable lift of B that is continuous in the second variable, but not necessarily in the first variable, such that $\Theta(X \times \{0\})$ is bounded. It follows that there exists a bounded measurable function $\psi : \mathcal{U} \times X \rightarrow \mathbb{Z}$ such that

$$F_{\Phi(\tilde{A}),2}(f(x), \Theta(x, t)) = \Theta(f(x), F_{\tilde{A},2}(x, t)) + \psi(\tilde{A}, x).$$

Thus $\rho'(\Phi(\tilde{A})) - \rho(A) = \int \psi(\tilde{A}, x) d\mu(x)$. On the other hand, $\psi(\tilde{A}, x)$ is easily seen to be continuous in the first variable, and since \mathcal{U} is connected, it depends only on the second variable. \square

APPENDIX D. PROOF OF THE PROJECTION LEMMA

Let us now prove Lemma 1.3. In the case where $\text{tr } A$ does not vanish identically, this result is contained in Lemma 9 of [ABD]. We will explain the necessary changes to deal with the remaining case.

We wish to show that an arbitrary continuous $\text{SL}(2, \mathbb{R})$ perturbation \tilde{A} of the constant cocycle A is conjugate to a Schrödinger perturbation (with a conjugacy depending continuously on the perturbation, and equal to id at A). On the other hand, given any fixed compact set of non-empty interior $K \subset X$, Lemma 10 of [ABD] states that under the assumption of minimality, \tilde{A} is conjugate to a perturbation \hat{A} such that $\hat{A} = A$ outside K .

So it is enough to describe how to conjugate such a localized perturbation. Since X has at least four points, we can choose K admitting a compact neighborhood K' such that $f^i(K') \cap K' = \emptyset$ for $1 \leq i \leq 3$. Notice that $A^4 = \text{id}$, so \hat{A}^4 is close to id . Define $E = \sup_{x \in X} \|\hat{A}^4(x) - \text{id}\|^{1/2}$.

For each $x \in K'$, there are unique $E_i(x) \in \mathbb{R}$, $1 \leq i \leq 4$, such that $E_3(x) = E\phi(x)$ and $S_{E_4(x)} \cdots S_{E_1(x)} = \hat{A}^4(x)$, where for $t \in \mathbb{R}$, S_t denotes the Schrödinger matrix $\begin{pmatrix} t & -1 \\ 1 & 0 \end{pmatrix}$. Indeed, write $\hat{A}^4(x) = R_{-\pi/2} \hat{A}(x) = \text{id} + \begin{pmatrix} a(x) & b(x) \\ c(x) & d(x) \end{pmatrix}$, and notice that $E_1(x)$, $E_2(x)$ and $E_4(x)$ must be given by

$$E_1(x) = \frac{b(x) - E\phi(x)}{1 + d(x)}, \quad E_2(x) = -\frac{d(x)}{E\phi(x)}, \quad E_4(x) = \frac{c(x) + E_2(x)}{1 - d(x)},$$

with the expression for $E_2(x)$ being evaluated as 0 whenever $E\phi(x) = 0$. These expressions immediately imply that the maps $E_i : K' \rightarrow \mathbb{R}$ are continuous and vanish in $\partial K'$.

We now define $\Phi(\hat{A})$ as follows. Outside $\bigcup_{i=0}^3 f^i(K')$, we set $\Phi(\hat{A}) = \hat{A} = A$. For $x \in K'$, we define $\Phi(\hat{A})(f^i(x)) = S_{E_{i+1}(x)}$, $0 \leq i \leq 3$. It is easy to see that Φ depends continuously on \hat{A} in a neighborhood of A .

By construction, $\Phi(\hat{A})^4(x) = \hat{A}^4(x)$ for $x \in K$. Define $\Psi(\hat{A}) = \text{id}$ outside $\bigcup_{i=0}^3 f^i(K')$, $\Psi(\hat{A})(f^i(x)) = \Phi(\hat{A})^i(x) \hat{A}^i(x)^{-1}$ for $x \in K'$ and $1 \leq i \leq 3$. Then Ψ is a continuous function of \hat{A} with $\Psi(\hat{A})(f(x)) \cdot \hat{A}(x) \cdot (\Psi(\hat{A})(x))^{-1} = \Phi(\hat{A})(x)$ for every $x \in X$. This completes the proof of Lemma 1.3.

Remark D.1. Like [ABD, Lemma 9], Lemma 1.3 is also valid in smooth settings. We only need to be careful in two steps. First, we should choose ϕ to be a smooth function. Second, we should replace the definition of $E = E(\hat{A})$, which should be “much larger than d ” for our argument to work, by something that takes into account the derivatives of d . In the case of C^r cocycles with finite r , we could take $E = \|d\|_{C^r}^{1/2}$, and it is straightforward to check that both Φ and Ψ are C^r -valued continuous maps defined in a C^r -neighborhood of A . In the case $r = \infty$, we can take

$$E = -\frac{1}{\log \sup_{x \in K'} \|\hat{A}^4(x) - \text{id}\|}$$

whenever $\hat{A} \neq A$, and let $E = 0$ when $\hat{A} = A$. With this definition, Φ and Ψ are C^0 -valued continuous maps in a C^0 -neighborhood of A , but their restrictions to C^{s+1} are C^s -valued continuous maps for every finite s (using the convexity inequality $\|\cdot\|_{C^s} \leq K_s (\|\cdot\|_{C^0} \|\cdot\|_{C^{s+1}}^s)^{1/(s+1)}$; see [N]), so their restrictions to C^∞ are C^∞ -valued continuous maps.

APPENDIX E. MIXING DYNAMICS AND CANTOR SPECTRUM

In [Fa], Fayad proved that there exists a real-analytic *time change* f^t of an irrational flow on $\mathbb{R}^3/\mathbb{Z}^3$ which is mixing. Thus f^t is generated by a vector field of the form $u\alpha$, where the coordinates of $\alpha \in \mathbb{R}^3$ are rationally independent and $u : \mathbb{R}^3/\mathbb{Z}^3 \rightarrow \mathbb{R}$ is a positive real analytic function. By general properties of time changes (see [Pa, §5.1]), to any measure μ_0 which is invariant by the flow generated by α , we can associate an f^t -invariant measure which is absolutely continuous with respect to μ_0 and has density $1/u$, and all f^t -invariant measures are obtained in this way. Since irrational flows are uniquely ergodic, the flow f^t is uniquely ergodic as well. We will denote by μ the unique f^t -invariant probability measure.

Proposition E.1. *Let $f = f^1$ be the time one map of the flow f^t . Then f is mixing, strictly ergodic, and the range G of the Schwartzman asymptotic cycle of f is dense in \mathbb{R} .*

Proof. Obviously f is mixing. If μ' is any invariant probability measure for f , then $\int_0^1 f_*^t \mu' dt$ is an f^t -invariant probability measure, so it coincides with μ . On the other hand, for every $\varepsilon > 0$, $\mu'_\varepsilon = \frac{1}{\varepsilon} \int_0^\varepsilon f_*^t \mu' dt$ is an f -invariant probability measure absolutely continuous with respect to μ . Obviously $\mu'_\varepsilon \rightarrow \mu'$ as $\varepsilon \rightarrow 0$, and since f is mixing, $f_*^n \mu'_\varepsilon \rightarrow \mu$ as $n \rightarrow \infty$ (in both cases the limits are taken in the weak-* topology). Since μ'_ε is f -invariant, we conclude that $\mu'_\varepsilon = \mu$ for every $\varepsilon > 0$, and hence $\mu' = \mu$. Thus f is uniquely ergodic. Since $\text{supp } \mu = \mathbb{R}^3/\mathbb{Z}^3$, unique ergodicity implies strict ergodicity.

Given $k \in \mathbb{Z}^3$, let $\psi : \mathbb{R}^3/\mathbb{Z}^3 \rightarrow \mathbb{R}/\mathbb{Z}$ be given by $\psi(x) = \langle x, k \rangle$, and let $\phi(x) = \int_0^1 u(f^t(x)) dt \langle \alpha, k \rangle$. Then $\phi(x) = \psi(f(x)) - \psi(x) \bmod 1$, $g = \int_{\mathbb{R}^3/\mathbb{Z}^3} \phi d\mu \in G$. Recalling that μ is an f^t -invariant probability measure with density proportional to $\frac{1}{u(x)}$, we get

$$\begin{aligned} \int_{\mathbb{R}^3/\mathbb{Z}^3} \phi(x) d\mu(x) &= \int_0^1 \int_{\mathbb{R}^3/\mathbb{Z}^3} u(f^t(x)) d\mu(x) dt \langle \alpha, k \rangle = \int_{\mathbb{R}^3/\mathbb{Z}^3} u(x) d\mu(x) \langle \alpha, k \rangle \\ &= C^{-1} \langle \alpha, k \rangle, \end{aligned}$$

where $C = \int_{\mathbb{R}^3/\mathbb{Z}^3} u(x)^{-1} dx > 0$ is the normalization constant making μ a probability measure. Since the coordinates of α are rationally independent, this shows that G contains a free Abelian subgroup of rank 3, so G is dense in \mathbb{R} . \square

It follows from the proposition and Corollary 2 that Schrödinger operators arising from mixing strictly ergodic dynamics can have Cantor spectrum, thus answering negatively [Be1, Problem 5].

REFERENCES

- [An] H. ANZAI. Ergodic skew product transformations on the torus. *Osaka Math. J.*, 3 (1951), 83–99.
- [At] K. ATHANASSOPOULOS. Some aspects of the theory of asymptotic cycles. *Exposition. Math.*, 13 (1995), 321–336.
- [AB] A. AVILA, J. BOCHI. A uniform dichotomy for generic $\text{SL}(2, \mathbb{R})$ cocycles over a minimal base. *Bull. Soc. Math. France*, 135 (2007), 407–417.

- [ABD] A. AVILA, J. BOCHI, D. DAMANIK. Cantor spectrum for Schrödinger operators with potentials arising from generalized skew-shifts. *Duke Math. J.*, 146 (2009), 253–280.
- [AF] A. AVILA, G. FORNI. Weak mixing for interval exchange transformations and translation flows. *Ann. of Math.*, 165 (2007), 637–664.
- [AJ1] A. AVILA, S. JITOMIRSKAYA. The Ten Martini Problem. To appear in *Ann. of Math.*
- [AJ2] A. AVILA, S. JITOMIRSKAYA. Almost localization and almost reducibility. To appear in *J. Eur. Math. Soc.*
- [AK] A. AVILA, R. KRIKORIAN. Reducibility or nonuniform hyperbolicity for quasiperiodic Schrödinger cocycles. *Ann. of Math.*, 164 (2006), 911–940.
- [AS1] J. AVRON, B. SIMON. Almost periodic Schrödinger operators. I. Limit periodic potentials. *Comm. Math. Phys.*, 82 (1981), 101–120.
- [AS2] J. AVRON, B. SIMON. Almost periodic Schrödinger operators. II. The integrated density of states. *Duke Math. J.*, 50 (1983), 369–391.
- [BCL] F. BÉGUIN, S. CROVISIER, F. LE ROUX. Construction of curious minimal uniquely ergodic homeomorphisms on manifolds: the Denjoy-Rees technique. *Ann. Sci. École Norm. Sup.*, 40 (2007), 251–308.
- [Be1] J. BELLISSARD. Schrödinger operators with almost periodic potential: An overview. *Mathematical Problems in Theoretical Physics (Berlin, 1981)*, 356–363, Lecture Notes in Phys., 153, Springer, Berlin-Heidelberg-New York, 1982.
- [Be2] J. BELLISSARD. K -theory of C^* -algebras in solid state physics. *Statistical Mechanics and Field Theory: Mathematical Aspects (Groningen, 1985)*, 99–156, Lecture Notes in Phys., 257, Springer, Berlin, 1986.
- [Be3] J. BELLISSARD. Gap labelling theorems for Schrödinger operators. *From Number Theory to Physics (Les Houches, 1989)*, 538–630, Springer, Berlin, 1992.
- [Be4] J. BELLISSARD. The noncommutative geometry of aperiodic solids. *Geometric and Topological Methods for Quantum Field Theory (Villa de Leyva, 2001)*, 86–156, World Sci. Publ., River Edge, NJ, 2003.
- [BBG] J. BELLISSARD, R. BENEDETTI, J.-M. GAMBAUDO. Spaces of tilings, finite telescopic approximations and gap-labeling. *Comm. Math. Phys.*, 261 (2006), 1–41.
- [BLT] J. BELLISSARD, R. LIMA, D. TESTARD. Almost periodic Schrödinger operators. *Mathematics + Physics. Vol. 1*, 1–64, World Sci. Publishing, Singapore, 1985.
- [BSc] J. BELLISSARD, E. SCOPPOLA. The density of states for almost periodic Schrödinger operators and the frequency module: a counterexample. *Comm. Math. Phys.*, 85 (1982), 301–308.
- [BSi] J. BELLISSARD, B. SIMON. Cantor spectrum for the almost Mathieu equation. *J. Funct. Anal.*, 48 (1982), 408–419.
- [BO] M.-T. BENAMEUR, H. OYONO-OYONO. Gap-labelling for quasi-crystals (proving a conjecture by J. Bellissard). *Operator Algebras and Mathematical Physics (Constanța, 2001)*, 11–22, Theta, Bucharest, 2003.
- [BJ] K. BJERKLÖV, T. JÄGER. Rotation numbers for quasiperiodically forced circle maps – mode-locking vs. strict monotonicity. *J. Amer. Math. Soc.*, 22 (2009), 353–362.
- [Bo] J. BOCHI. Genericity of zero Lyapunov exponents. *Ergod. Th. Dynam. Sys.*, 22 (2002), 1667–1696.
- [BC] C. BONATTI, S. CROVISIER. Récurrence et générnicité. *Invent. Math.*, 158 (2004), 33–104.
- [Ca] R. H. CAMERON. Linear differential equations with almost periodic coefficients. *Ann. of Math.*, 37 (1936), 29–42.
- [CEY] M. D. CHOI, G. A. ELLIOTT, N. YUI. Gauss polynomials and the rotation algebra. *Invent. Math.*, 99 (1990), 225–246.
- [Ch] V. A. CHULAEVSKIĬ. An inverse spectral problem for limit-periodic Schrödinger operators. *Funktional. Anal. i Prilozhen.*, 18 (1984), 63–66.
- [CS] V. A. CHULAEVSKIĬ, YA G. SINAĬ. Anderson localization for the 1-D discrete Schrödinger operator with two-frequency potential. *Comm. Math. Phys.*, 125 (1989), 91–112.
- [CF] N. D. CONG, R. FABBRI. On the spectrum of the one-dimensional Schrödinger operator. *Discrete Contin. Dyn. Syst.*, 9 (2008), 541–554.
- [Co] A. CONNES. An analogue of the Thom isomorphism for crossed products of a C^* -algebra by an action of \mathbb{R} . *Adv. in Math.*, 39 (1981), 31–55.
- [DS] F. DELYON, B. SOUILLARD. The rotation number for finite difference operators and its properties. *Comm. Math. Phys.*, 89 (1983), 415–426.
- [E] L. H. ELIASSON. Floquet solutions for the 1-dimensional quasi-periodic Schrödinger equation. *Comm. Math. Phys.*, 146 (1992), 447–482.
- [EJ] R. ELLIS, R. JOHNSON. Topological dynamics and linear differential systems. *J. Differential Equations*, 44 (1982), 21–39.

- [FJ] R. FABBRI, R. JOHNSON. Genericity of exponential dichotomy for two-dimensional differential systems. *Ann. Mat. Pura Appl.*, 178 (2000), 175–193.
- [FJP] R. FABBRI, R. JOHNSON, R. PAVANI. On the nature of the spectrum of the quasi-periodic Schrödinger operator. *Nonlinear Anal. Real World Appl.*, 3 (2002), 37–59.
- [Fa] B. FAYAD. Analytic mixing reparametrizations of irrational flows. *Ergodic Theory Dynam. Systems*, 22 (2002), 437–468.
- [Fu] H. FURSTENBERG Strict ergodicity and transformation of the torus. *Amer. J. Math.*, 83 (1961), 573–601.
- [GS] M. GOLDSTEIN, W. SCHLAG. On resonances and the formation of gaps in the spectrum of quasi-periodic Schrödinger equations. Preprint <http://arxiv.org/abs/math/0511392>.
- [HS] B. HÉLFFER, J. SJÖSTRAND. Semiclassical analysis for Harper’s equation. III. Cantor structure of the spectrum. *Mém. Soc. Math. France*, 39 (1989), 1–124.
- [H] M.-R. HERMAN. Une méthode pour minorer les exposants de Lyapounov et quelques exemples montrant le caractère local d’un théorème d’Arnol’d et de Moser sur le tore de dimension 2. *Comment. Math. Helv.*, 58 (1983), 453–502.
- [J] R. JOHNSON. Exponential dichotomy, rotation number, and linear differential operators with bounded coefficients. *J. Differential Equations*, 61 (1986), 54–78.
- [JM] R. JOHNSON, J. MOSER. The rotation number for almost periodic potentials. *Comm. Math. Phys.*, 84 (1982), 403–438.
- [KP] J. KAMINKER, I. PUTNAM. A proof of the gap labeling conjecture. *Michigan Math. J.*, 51 (2003), 537–546.
- [KZ] J. KELLENDONK, I. ZOIS. Rotation numbers, boundary forces and gap labelling. *J. Phys. A*, 38 (2005), 3937–3946.
- [K] O. KNILL. The upper Lyapunov exponent of $SL(2, \mathbb{R})$ cocycles: discontinuity and the problem of positivity. *Lyapunov Exponents (Oberwolfach, 1990)*, 86–97, Lecture Notes in Math., 1486, Springer, Berlin, 1991.
- [L] Y. LAST. Zero measure spectrum for the almost Mathieu operator. *Comm. Math. Phys.*, 164 (1994), 421–432.
- [MMY] S. MARMI, P. MOUSSA, J.-C. Yoccoz. The cohomological equation for Roth-type interval exchange maps. *J. Amer. Math. Soc.*, 18 (2005), 823–872.
- [M] J. MOSER. An example of a Schrödinger equation with almost periodic potential and nowhere dense spectrum. *Comment. Math. Helv.*, 56 (1981), 198–224.
- [N] L. NIRENBERG. On elliptic partial differential equations. *Ann. Scuola Norm. Sup. Pisa*, 13 (1959), 115–162.
- [PT] L. A. PASTUR, V. A. TKACHENKO. Spectral theory of a class of one-dimensional Schrödinger operators with limit-periodic potentials. *Trans. Moscow Math. Soc.*, 1989 (1989), 115–166.
- [PV] M. PIMSNER, D. VOICULESCU. Exact sequences for K -groups and Ext-groups of certain cross-product C^* -algebras. *J. Operator Theory*, 4 (1980), 93–118.
- [Pa] W. PARRY. *Topics in Ergodic Theory*. Cambridge University Press, 1981.
- [Pu] J. PUIG. Cantor spectrum for the almost Mathieu operator. *Comm. Math. Phys.*, 244 (2004), 297–309.
- [Sc] S. SCHWARTZMAN. Asymptotic cycles. *Ann. of Math.*, 66 (1957), 270–284.
- [Si1] B. SIMON. On the genericity of nonvanishing instability intervals in Hill’s equation. *Ann. Inst. H. Poincaré*, A24 (1976), 91–93.
- [Si2] B. SIMON. Almost periodic Schrödinger operators: A review. *Adv. Appl. Math.*, 3 (1982), 463–490.
- [X] J. XIA. The Floquet theory and the state density of quasi-periodic Schrödinger operators. *J. Differential Equations*, 86 (1990), 243–259.
- [Y] J.-C. YOCOZO. Some questions and remarks about $SL(2, \mathbb{R})$ cocycles. In *Modern Dynamical Systems and Applications*, 447–458. Cambridge Univ. Press, Cambridge, 2004.

URL: www.impa.br/~avila

E-mail address: artur@math.sunysb.edu

URL: www.mat.puc-rio.br/~jairo

E-mail address: jairo@mat.puc-rio.br

URL: www.ruf.rice.edu/~dtd3

E-mail address: damanik@rice.edu